Assignment 2: CS 215, Fall 2024

Prof. Sunita Sarawagi

Released: 31 August, 2024

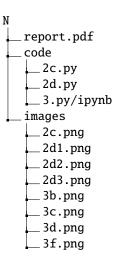
Hints and Soft Deadline: 9 September, 2024 Hard Deadline: 11:59 PM, 13 September, 2024

PLEASE READ ALL INSTRUCTIONS CAREFULLY!

Instructions

- 1. **Report instructions**. You are requested to type out your solutions in LaTeX. If this is not possible, you may print out a2_hand.pdf, hand-write all solutions neatly in the space provided and submit a pdf containing scanned copies of the filled-in document. Again, please note that humans aren't perfect image-to-text converters; do write large and clear!
- 2. **Report instructions, continued**. The submitted file must contain the names and roll numbers of all group members on the first page. Some problems will require code; please make sure adequate running instructions are included with your solution to each such problem.
- 3. **Code file instructions**. For each code file: type out the name and roll number of each group member in a comment at the top of the file, along with the problem number the code is for.
- 4. And then, check format. Once the pdf and code are ready, zip them into one folder with the name N=A2-RollNumberOfFirstStudent-RollNumberOfSecondStudent-RollNumberofThirdStudent.

 Please name the zip file N.zip. You may use the command zip -r N.zip N from the directory containing your submission directory N to do this. If the assignment was done in a group of size not equal to 3, the name of the folder and zip file will be A2 followed by the roll numbers of each member; each item separated by hyphens. Any letters in the roll number must be capitalized. The submissions for this assignment MUST follow the given directory structure and assignments not adhering to this risk not being graded.



- 5. and submit (on time!). Upload the file on moodle BEFORE 11:59 pm on the due date. No assignments will be accepted thereafter. Please expect Murphy's Law to hold post 11.45 pm on the due date. Note that only one student per group should upload their work on moodle, though all group members will receive grades.
- 6. **Misc 1: Bonuses**. Some tasks in some questions are optional and are marked with a **B**. Solving some of these questions can get you extra points, capped to the maximum number of points on the assignment.

- 7. **Misc 2: Hints**. After a little over a week, hints will be posted to some tasks of some questions (precisely the tasks marked with a ★). All submissions that are not updated on moodle at any point after the soft deadline will receive 10 credits of bonus points (capped off at the maximum number of points for the assignment), for not using the hints. Use the trade-off of scoring 10 points without hints and perhaps scoring more than that with a few well-placed hints wisely!
- 8. **Total**. The maximum possible marks on this assignment are 100 with the marks being normally distributed between questions, that is, Q1 and Q5 carry 15 marks each, Q2 and Q4 carry 20 marks each and Q3 carries 30 marks. Let us know what you think the "most probable" mean and variance is for this Gaussian "histogram" ;-P
- 9. **Plagiarism**. This assignment is mostly for your own learning, and you are free to discuss the problems with anyone. However, verbatim copying of solutions is unacceptable: compromising your credibility for 5-10 points in one course (whose grade will not even matter in a few years from now) is not worth it by a mile. Please give the problems your time; they are meant to help you understand the subject better.
- 10. One last thing. Please preserve a copy of all your work until the end of the semester. And please have fun!

MATHEMAGIC 3

§ 1 Mathemagic

We explore some connections between some special random variables, via the notion of probability generating functions.

Definition 1 (PGF, MGF). Let X be a random variable taking on only non-negative-integer values. Suppose that X is distributed according to probability mass function P. We define the probability-generating-function of the distribution P[X] of random variable X by

$$G(z) := \mathbb{E}\left[z^X\right] = \sum_{n=0}^{\infty} P\left[X = n\right] z^n.$$

The moment-generating-function of the same distribution is defined by $M(t) := G(e^t)$ for t such that the series for G converges.

♦ Task A [1]

Derive the PGF when X is a Bernoulli random variable with parameter p, that is, $X \sim \text{Ber}(p)$. Call this PGF G_{Ber} .

Solution

$$G_{\text{Ber}}(z) = (1-p)z^0 + pz^1 = 1-p+pz.$$

Rubric: 1 point for the correct answer.

 \Diamond Task B [2]

Let G_{Bin} the PGF when $X \sim Bin(n, p)$. Show that $G_{Bin}(z) = G_{Ber}(z)^n$.

Solution

$$G_{\text{Bin}}(z) = \sum_{r=0}^{n} P[X = r] z^{r} = \sum_{r=0}^{n} {n \choose r} (1-p)^{n-r} (pz)^{r} = (1-p+pz)^{n} = G_{\text{Ber}}(z)^{n}.$$

Rubric: 1 point for the correct answer, 1 point for the correct derivation.

 $\diamond \operatorname{Task} \mathsf{C} (\star)$ [4]

Suppose that X_1, X_2, \dots, X_k are independent non-negative-integer-valued random variables, each distributed with the same probability mass function P. Let G be their common PGF. Consider random variable $X = X_1 + X_2 + \dots + X_k$ defined on the cartesian product of the sample spaces underlying the random variables X_i . Let the PGF corresponding to X be G_{Σ} . Show that $G_{\Sigma}(z) = G(z)^k$.

Solution

Suppose $X_i : \Omega_i \to \mathbb{R}$ are defined on sample spaces Ω_i . Then $X : \Omega_1 \times \Omega_2 \times \cdots \times \Omega_k \to \mathbb{R}$ is defined on the cartesian product by $X(s_1, \dots, s_k) := X_1(s_1) + X_2(s_2) + \cdots + X_k(s_k)$.

Solution 1. Let $G(z) = \sum_{i=0}^{\infty} p_i z^i$. Then the coefficient of z^n in $G(z)^k$ is

$$[z^n]G(z)^k = \sum_{\substack{i_1,\dots,i_k\\i_1+i_2+\dots+i_k=n}} p_{i_1}p_{i_2}\cdots p_{i_k} = \sum_{\substack{i_1,\dots,i_k\\i_1+i_2+\dots+i_k=n}} P[X_1=i_1,X_2=i_2,\dots,X_k=i_k].$$

But the last expression is precisely the probability that X = n. It follows that for each $n \ge 0$, $P[X = n] = [z^n]G(z)^k$, i.e. $G_{\Sigma}(z)$ and $G(z)^k$ agree coefficient-wise. So $G_{\Sigma}(z) = G(z)^k$.

MATHEMAGIC 4

Solution 2. We prove the theorem for the case k = 2; the general case follows by induction on k, noting that the variables $X_1 + X_2 + \cdots + X_{k-1}$ and X_k are independent. For the case k = 2, we have

$$G_{\Sigma}(z) = \sum_{n=0}^{\infty} P[X = n] z^n = \sum_{n=0}^{\infty} \sum_{i=0}^{n} P[X_1 = i, X_2 = n - i] z^n = \sum_{i=0}^{\infty} \sum_{j=0}^{\infty} P[X_1 = i] z^j P[X_2 = j] z^j = G(z)G(z).$$

Solution 3. We note that $\mathbb{E}[XY] = \mathbb{E}[X]\mathbb{E}[Y]$ for independent random variables X, Y, since

$$\mathbb{E}\left[XY\right] = \sum_{x,y} xy P\left[X = x, Y = y\right] = \left(\sum_{x} x P\left[X = x\right]\right) \left(\sum_{y} y P\left[Y = y\right]\right) = \mathbb{E}\left[X\right] \mathbb{E}\left[Y\right].$$

Also, it can be verified that the formal random variables z^X and z^Y are independent for independent random variables X, Y. Thus, by the expectation-value definition of the PGF,

$$G_{\Sigma}(z) = \mathbb{E}\left[z^{X}\right] = \mathbb{E}\left[z^{X_{1} + X_{2} + \dots + X_{k}}\right] = \mathbb{E}\left[z^{X_{1}}\right] \mathbb{E}\left[z^{X_{2}}\right] \cdots \mathbb{E}\left[z^{X_{k}}\right] = G(z)^{k}.$$

Rubric: 4 points for any correct solution. Partial credit for partial progress.

♦ Task D [1]

Consider now the Geometric distribution. Let $X \sim \text{Geo}(p)$. Derive its PGF.

Solution

$$G_{\text{Geo}}(z) = \sum_{n=1}^{\infty} P[X = n] z^n = \sum_{n=1}^{\infty} (1-p)^{n-1} p z^n = p z \sum_{n=0}^{\infty} ((1-p)z)^n = \frac{pz}{1-(1-p)z}.$$

Rubric: 1 point for the correct answer.

♦ Task E [3]

Consider $X \sim \text{Bin}(n, p)$ and $Y \sim \text{NegBin}(n, p)$. Let their PGFs be $G_X^{(n,p)}(z)$ and $G_Y^{(n,p)}$ respectively. Show using previous tasks or otherwise that for every 0 , we have

$$G_Y^{(n,p)}(z) = \left(G_X^{(n,p^{-1})}(z^{-1})\right)^{-1}.$$

Solution

We can write $Y = X_1 + X_2 + \cdots + X_n$ where each $X_i \sim \text{Geo}(p)$ are independent. Then by Task C,

$$G_Y^{(n,p)}(z) = G_{\text{Geo}}(z)^n = \left(\frac{pz}{1-(1-p)z}\right)^n = \left(\frac{1}{1-p^{-1}+p^{-1}z^{-1}}\right)^n = \left(G_X^{(n,p^{-1})}(z^{-1})\right)^{-1}.$$

Rubric: 3 points for the correct derivation.

That shows that the negative binomial distribution is not only morally an "inverse" of the binomial distribution (in that it models number of trials while fixing number of successes, while the binomial fixes the number of trials and models the number of successes), but negates the binomial distribution in every way possible! There is more, see Task F below.

♦ Task F [2]

We shall derive the negative binomial theorem using the result from Task E. First, we generalize the binomial coefficient: let $\alpha \in \mathbb{R}$ and $k \in \mathbb{N}$. Then

$$\binom{\alpha}{k} := \frac{\alpha(\alpha-1)\cdots(\alpha-k+1)}{k!}.$$

MATHEMAGIC 5

Theorem 2 (Binomial theorem, negative exponent). *Let* $n \in \mathbb{N}$ *and* |x| < 1. *Then*

$$(1+x)^{-n} = \sum_{r=0}^{\infty} (-1)^r \binom{n+r-1}{r} x^r = \sum_{r=0}^{\infty} \binom{-n}{r} x^r.$$

Solution

The key idea is that we will compute the PGF for the negative binomial distribution another way, and equate it with the result from Task E. Suppose that $Y \sim \text{NegBin}(n, p)$. Then

$$P[Y = r] = \binom{r-1}{n-1} p^n (1-p)^{r-n}.$$

So the PGF is

$$G_Y(z) = \sum_{r=n}^{\infty} \binom{r-1}{n-1} p^n (1-p)^{r-n} z^r = (pz)^n \sum_{r=0}^{\infty} \binom{r+n-1}{n-1} ((1-p)z)^r.$$

However, from Task E, the PGF is also

$$G_Y(z) = \left(G_X^{(n,p^{-1})}(z^{-1})\right)^{-1} = \left(\frac{1}{1-p^{-1}+p^{-1}z^{-1}}\right)^n = (pz)^n \left(1-(1-p)z\right)^{-n}.$$

Letting x = -(1 - p)z, we have by equating both expressions that

$$(1+x)^{-n} = \sum_{r=0}^{\infty} \binom{r+n-1}{n-1} (-x)^r.$$

The final step is to note that

$$\binom{-n}{r} = \frac{(-n)(-n-1)\cdots(-n+r-2)(-n+r-1)}{r!} = (-1)^r \frac{(n+r-1)(n+r-2)\cdots(n+1)n}{r!} = (-1)^r \binom{n+r-1}{r},$$

from which the result immediately follows.

Digression: this kind of argument, where one quantity is computed in two different ways, often leads to new proofs of well-known results. In many cases, however, the "new proof" turns out to be essentially a rephrasing of an earlier proof. Is this proof equivalent to a proof some of you have seen before of the negative binomial theorem?

Rubric: 2 points for the correct derivation.

It is almost (the summation not being finite is the only difference) as if one could simply put a negative number for the exponent in the usual binomial theorem.

$$\diamond$$
 Task G [2]

An easy consequence of all the hard work. Suppose that the PGF of random variable X is G(z). Show that the expectation of X is simply the derivative of G at 1, i.e. $\mathbb{E}[X] = G'(1)$. Using this and the PGFs constructed previously, derive the means of the Bernoulli, Binomial, Geometric and Negative binomial distributions as a function of their parameters.

Solution

The derivative of PGF $G(z) = \sum_{i=0}^{\infty} p_i z^i$ is $G'(z) = \sum_{i=1}^{\infty} i p_i z^{i-1}$, which is $\sum_{i=0}^{\infty} i p_i = \mathbb{E}[X]$ at z=1, as needed. A fancier way of writing the same thing is $G'(z) = \left(\mathbb{E}\left[z^X\right]\right)' = \mathbb{E}\left[Xz^{X-1}\right]$, so $G'(1) = \mathbb{E}[X]$. The rest is just calculus, simply differentiate the PGFs treating them as functions of z (they are not really functions of z, but the rules of calculus nevertheless apply). The answers are

$$\mathbb{E}[X] = p$$
, $\mathbb{E}[X] = np$, $\mathbb{E}[X] = \frac{1}{p}$, $\mathbb{E}[X] = \frac{n}{p}$.

Rubric: 1/2 points per correct answer, if explanation provided for $\mathbb{E}[X] = G'(1)$.

§ 2 Normal Sampling

It's nice and all to know that standard normal random variables exist and are well-behaved; but can we actually sample from a standard normal distribution?

Let's be clear on what we want to do. We would like to find an algorithm \mathcal{A} that takes some uniformly random numbers in [0,1], and generates an output $x \in \mathbb{R}$. The algorithm should use the randomness of the numbers it receives in such a way that the output x of the algorithm is distributed standard normally. Denote by $f_{\mathcal{A}}(x)$ the PDF of the algorithm's output (a random variable). We wish that for every $x \in \mathbb{R}$,

$$f_{\mathcal{H}}(x) = \frac{1}{\sqrt{2\pi}}e^{-x^2/2}.$$

To find such an algorithm, we will show and then use a very general theorem. A note on notation: we shall use f_X to denote the PDF of random variable X, and F_X to denote the CDF.

♦ Task A [2]

Theorem 3. Let X be a continuous real-valued random variable with CDF $F_X : \mathbb{R} \to [0,1]$. Assume that F_X is invertible. Then the random variable $Y := F_X(X) \in [0,1]$ is uniformly distributed in [0,1].

Prove the theorem above.

Solution

Let F_Y denote the CDF of Y. Then for any $u \in [0,1]$,

$$F_Y(u) = {\rm P} \big[Y \leq u \big] = {\rm P} \big[F_X(X) \leq u \big] = {\rm P} \big[X \leq F_X^{-1}(u) \big] = F_X(F_X^{-1}(u)) = u,$$

thus $f_Y(u) = F_Y'(u) = 1$ for every $u \in [0, 1]$ or Y is uniformly distributed in [0, 1]. **Rubric:** 2 points for the correct proof.

 \Diamond Task B [2]

Suppose now that we are given random variable Y which is uniform over [0,1]. Explain how an algorithm \mathcal{A} may be constructed taking as input a sample y according to the distribution of Y (so the algorithm is given just one uniformly random number in [0,1]), such that for every $u \in \mathbb{R}$, we have

$$F_{\mathcal{A}}(u) = F_X(u)$$
.

In other words, the output of \mathcal{A} has the same cumulative distribution function as X, and upon taking the derivative so has the same PDF.

Solution

Consider algorithm \mathcal{A} that outputs $F_X^{-1}(y)$ upon input y. Treating the output of the algorithm as a random variable (also denoted \mathcal{A}), we have $\mathcal{A} = F_X^{-1}(Y)$. Then for any $u \in \mathbb{R}$,

$$F_{\mathcal{A}}(u) = \mathrm{P}\left[\mathcal{A} \leq u\right] = \mathrm{P}\left[F_X^{-1}(Y) \leq u\right] = \mathrm{P}\left[Y \leq F_X(u)\right] = F_Y(F_X(u)) = F_X(u),$$

where the last equality follows from the fact that Y is uniform on [0, 1].

Rubric: 2 points for the correct algorithm. Any other correct algorithm is also acceptable.

♦ Task C [8]

In this task, you will use the numpy.random module to generate random numbers uniformly in [0,1] and the norm class (gives you access to F_X and F_X^{-1} for a Gaussian random variable X) from scipy.stats to sample from a Gaussian. In particular,

- Write a Python function sample(loc, scale) that samples from the Gaussian with mean at x = loc and standard deviation scale.
- Generate $N=10^5$ independent samples using the function above for the four parameter choices $(\mu, \sigma^2) = (0,0.2), (0,1.0), (0,5.0), (-2,0.5)$.
- Plot the samples for each parameter choice using matplotlib.pyplot. You should roughly reproduce the shape of the plot of the four Gaussians from the lecture slides, see figure 1. The bottom plot is a plot of the samples; it mimics the top plot of the different PDFs, confirming that we have indeed sampled from the different Gaussian PDFs. Save the plot in 2c.png.

You may not use the built in numpy.random.normal or random.randn functions that directly sample from Gaussians for this task. For the plot of the samples, you may not use kernel-density estimation (kde) tools; matplotlib.pyplot by itself is sufficient. Please write all your code for this task in one file called 2c.py.

Solution

We simply run the algorithm $\mathcal A$ from Task B with F_X^{-1} being replaced by scipy.stats.norm.ppf.

```
import numpy as np
import scipy.stats.norm as norm

def sample(loc, scale, N):
    return norm.ppf(np.random.random(N), loc=loc, scale=scale)
```

The plotting is routine using matplotlib.pyplot.hist.

Rubric: 2 point for explaining the idea, 3 points for correct code (working and not using disallowed functions), 3 points for the correct plot.

 $\Diamond \operatorname{Task} D(\star)$ [8]

Now we consider another way to sample normal distributions - approximately - using a Galton board (see figure 2 for a picture). Imagine a ball that starts at the top of a Galton board. As it hits the first peg, it moves to the left of the peg with probability 1/2 and to the right of the peg with probability 1/2. After falling vertically downwards a little, it hits another peg. Again, it moves left or right of the new peg with equal probability. Suppose it makes h collisions with pegs before reaching the bottom of the Galton board. We call h the depth of the board. At the end, the ball falls into one of the wood-piece-separated pockets. There are h+1 pockets at the bottom of the board, and each ball falls into exactly one of these as per the random directions it chose at each collision.

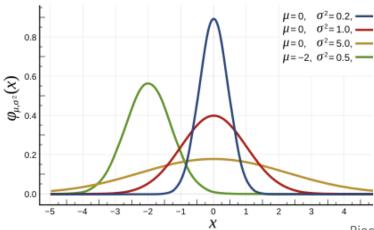
Consider simulating the motion of a large number N of balls along the Galton board, and record the fraction of balls that finally end up in each of the h + 1 pockets.

The simulation works by simulating the motion of each ball from the top to bottom of the Galton board. Initially the ball is at x = 0, and in the first step the ball moves left or right with equal probability - so its current position x is incremented or decremented with equal probability. The second step is identical - with a different starting position x. Suppose in the first step x was incremented, so x = 1 now. Then on the second step, x is decremented (back to 0) or incremented (to 2) with equal probability. Perhaps it was incremented again. The third step decrements it to 1 or increments it to 3 with equal probability. Perhaps it was decremented to 1. A fourth step takes it to 0 or 2 with equal probability. And so on. h steps ensue till the final value of its position x is the pocket the ball will fall into.

The choice of moving left or right can be simulated by simulating from $\{0,1\}$ uniformly randomly (find a function in numpy.random to do this). h uniform samples from $\{0,1\}$ thus allow us to simulate one ball.

The process can be repeated *N* times with the final pockets of each simulated ball being recorded, to yield a simulation of the Galton board. Note that it is possible, using numpy, to omit any Python loops. You are not required to omit loops in your code but it needs to be fast enough to gather the data required.

Your task is to carry out this simulation for $N = 10^5$ (if this is too large, you may use $N = 10^3$) for the three values h = 10, 50, 100 of the depth of the board. For each value of h, plot the counts of balls in each pocket obtained as a histogram, with one bin for each pocket. The code is to be written in file 2d.py, with the three histograms



Piech & Cain, CS109, Stanford University

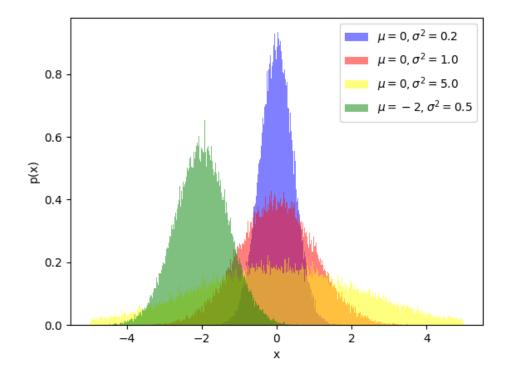


Figure 1: Question 2, Task C: The Bell curve, from uniformly random numbers

saved to files 2d1.png, 2d2.png and 2d3.png (see figure 2 for a reference plot). What do you notice about the shape of the tops of the histogram?

Solution

The key idea is that the final position of the ball is the sum of h independent random variables, each of which is -1 or 1 with equal probability. The code naturally follows.

```
# simulates N balls on a Galton board of depth h
def galton(h, N):
    return np.sum(np.random.choice([-1, 1], size=(h, N)), axis=0)
```

The plotting is routine. The important thing to note here is that not every integer value between -h and h is a pocket; indeed, integers having opposite parity to h will be empty, since no ball can reach these positions. The pockets are $\{-h, -h + 2, \cdots, 0, 2, \cdots, h - 2, h\}$, which must be the x axis when plotting the histogram.



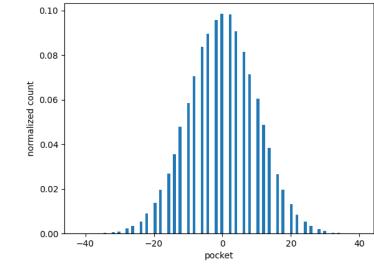


Figure 2: (top) A Galton board with h = 10 (ignore the leftmost and rightmost pockets). (bottom) A simulation with $N = 10^5$ and h = 100.

However, since this is a matter of style of plotting, there will be no penalty for plotting the histogram with all integers between -h and h on the x axis.

The shape of the tops of the histogram is an unmistakable Bell curve.

Rubric: 2 points for explaining the idea, 3 points for correct code (working and not using disallowed functions). Deduct one point if $N = 10^3$ is used instead of $N = 10^5$. 1 point for each correct histogram.

 $\Diamond \operatorname{Task} E (B)$ [5]

The goal here is to theoretically show that the distribution of number of balls in each pocket is normally distributed. Suppose that h = 2k is even. Consider a Galton board of depth h. Let random variable $X \in \{-h, -h + 2, \cdots, 0, 2, \cdots, h - 2, h\}$ describe the pocket in which a ball finally lands (notice that X can only be even). The probability mass distribution $P_h[\cdot]$ of X is determined by the randomness of the outcome of each collision. There are two sub-tasks here.

- Compute (in closed form) the value P_h [X = 2i] for each $i \in \{-k, -k+1, \dots, k-1, k\}$.
- Show that for large enough h and $i \ll \sqrt{h}$,

$$P_h[X = 2i] \approx \frac{1}{\sqrt{\pi k}} e^{-\frac{i^2}{k}} = \mathcal{N}(\mu = 0, \sigma^2 = k/2)(i).$$

You may need to use Stirling's approximation for the factorial

$$n! \approx \sqrt{2\pi n} \left(\frac{n}{e}\right)^n$$

to solve this subtask.

Solution

Suppose there were r rightward movements and l leftward movements leading to final position x = 2i. Then r + l = h and r - l = 2i. Solving these equations gives r = k + i and l = k - i. The probability of this happening is

$$P_h[X = 2i] = {h \choose r} \frac{1}{2^h} = \frac{{2k \choose k+i}}{2^{2k}}.$$

The second part is a bit more involved. We take (natural) logarithms on both sides.

$$\log P_h[X = 2i] = \log {2k \choose k+i} - 2k \log 2 = \log(2k)! - \log(k+i)! - \log(k-i)! - 2k \log 2.$$

Stirling's approximation can only become tighter when we take the logarithm, so we have

$$\log n! \approx \log \left(\sqrt{2\pi n} \left(\frac{n}{e} \right)^n \right) = \frac{1}{2} \log 2\pi + \frac{1}{2} \log n + n \log n - n.$$

Substituting this into the above equation and simplifying gives

$$\log P_h \left[X = 2i \right] \approx -\frac{1}{2} \log 2\pi + \frac{1}{2} \log \left(\frac{k^2 - i^2}{2k} \right) + (k+i) \log \frac{k}{k+i} + (k-i) \log \frac{k}{k-i}.$$

The first term is easy to handle; $i^2 \ll k$, i.e. $i^2 = o(k)$, so $(k^2 - i^2)/2k = k/2 + o(1)$. The logarithm is then $\log(k/2) + o(1)$. For the second and third terms, we use a Taylor approximation $\log(1 + x)$. We can write the latter two terms as

$$T = -(k+i)\log(1+i/k) - (k-i)\log(1-i/k).$$

Since $x = i/k = o(1/\sqrt{k})$ is small, we can use Taylor approximation. A common fallacy is to approximate $\log(1+x)$ by x, in which case T simplifies to $-2i^2/k$. However, this is erroneous; the approximation was correct up to $O(x^2) = O(i^2/k^2)$, so the error in T would be $O(i^2/k)$, which is comparable to the value of T itself - this is not a valid approximation. The correct approximation is to use $\log(1+x) \approx x - x^2/2$, which gives $T = -i^2/k + O(i^2/k^2)$. Here, the error term is $o(i^2/k)$, so we can ignore it. The final result is

$$\log P_h [X = 2i] \approx -\frac{1}{2} \log \pi k - \frac{i^2}{k},$$

exponentiating which yields the answer.

Rubric: 1 points for the correct derivation of $P_h[X=2i]$, 4 points for the correct derivation of the approximation. Partial credit for partial progress. The approximation is the key part of the solution, incorrect approximations receive only 1 point for the second part.

§ 3 Fitting Data

In this problem, we are given a dataset that came from some distribution that we don't know offhand, and we want to find what that distribution is. Why? So that we can then predict future outputs that are sampled from the same distribution, if we get close enough to the right answer.

Suppose we have a dataset consisting of n samples $S = \{x_1, \dots, x_n\}$, with each $x_i \in \mathbb{R}^d$. You are given a dataset which we will call \mathcal{D} in the file 3.data. We assume that there is an underlying probability distribution function P for a random variable X such that each x_i was sampled independently from P[X]: the i^{th} sample is a random variable X_i with the same distribution as X, i.e. $P[X_i = x_i] = P[X = x_i]$. Notice that a dataset may then simply be treated as the collection of outcomes of a number n of experiments, each of which consist of sampling X from its underlying distribution P.

Guessing or estimating the distribution P will tell us roughly what X_{n+1}, X_{n+2}, \cdots might turn out to be. This is useful, and so trying to estimate an unknown distribution from a sample will be our goal for this question.

Code for all subparts is to be written in 3.py. If you wish, you may choose to display plots with the code, using an .ipynb file. Please name it 3.ipynb in this case.

$$\Diamond \operatorname{Task} A$$
 [2]

The *i*th moment of a random variable is defined by $\mu_i := \mathbb{E}\left[X^i\right]$. When given a sample of n datapoints, we define the sample *i*th moment as the moment of the random variable that is each datapoint with equal probability 1/n. In particular, the sample *i*th moment is

$$\hat{\mu}_i := \frac{x_1^i + \dots + x_n^i}{n}.$$

Compute the first two (i = 1, 2) moments of \mathcal{D} . You may use numpy arrays and operations on them, but no loops are allowed to compute the moments (yes, it can be done with numpy).

♦ Task B [2]

Let's first try to graphically guess the mode of the distribution from our dataset.

Recall that a dataset is simply the collection of outcomes of many independent identical experiments, and that probability of an element x models the fraction of experiments with outcome x.

Compute a histogram of the dataset and plot it. From the histogram, graphically guess the mode(s) from the data. Please attach any code used in file 3.py, saving the histogram computed in file 3b.png.

A histogram is typically used to "see" what kind of distribution the sample could have come from.

Solution

The mode of the actual distribution is, likely, the peak of the histogram of the samples - in this case, around 6. The histogram code is simply plt.hist(data) with a suitable number of bins.

Rubric: 1 point for reasonably correct mode(s), 1 point for the correct histogram.

 \diamond Task C (\star) [1+2+2]

We have still got nowhere close to finding a distribution that fits the data we have well. Let's guess. It looks reasonably like a binomial distribution centred near about 6. Let us find the binomial distribution that is closest to our data.

One way to define closeness is to ask for the first few moments to be equal (it is a theorem in statistics that if every moment of two distributions are identical, then the two distributions must be identical as well - we are roughly approximating this theorem here).

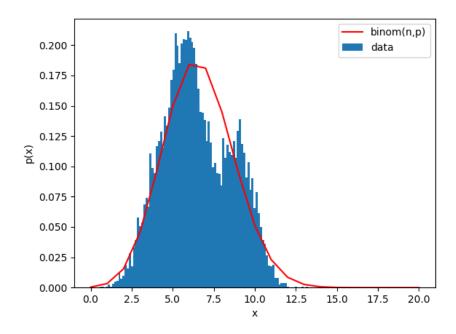


Figure 3: The best binomial distribution approximation to the true distribution

In this task, we will find the best-fit binomial distribution to \mathcal{D} by asking for the first two moments $\hat{\mu}_1$ and $\hat{\mu}_2$ to be equal to the corresponding two moments of Bin(n, p), for some suitable choice of n and p.

- 1. First, compute an expression for the first two moments μ_1^{Bin} , μ_2^{Bin} of the distribution Bin(n,p) as a function of n and p.
- 2. Then, use the fsolve function from scipy.optimize to compute a solution (n, p) to $\hat{\mu}_i = \mu_i^{\text{Bin}}$, i = 1, 2. Round n to either $\lfloor n \rfloor$ or $\lceil n \rceil$ based on which one satisfies the equalities better (they will not be exactly satisfied). Say the found parameters are (n^*, p^*) .
- 3. Finally, using numpy.linspace and scipy.stats.binom.pmf, plot the binomial distribution $Bin(n^*, p^*)$ on top of the histogram of \mathcal{D} . You should get something like the figure in figure 3. Save the plot to 3c.png.

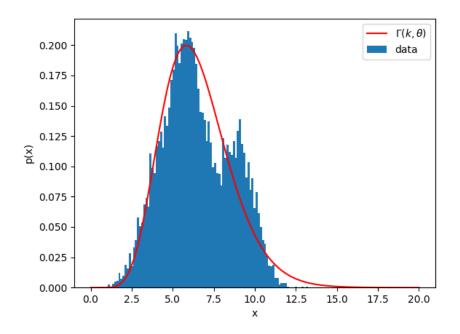


Figure 4: The best Gamma distribution approximation to the true distribution

 \diamond Task D [3+2+2]

Well, that was not too bad an approximation. Let us try another distribution, this time continuous.

The gamma distribution is a two-parameter family of continuous probability distributions. It is parameterized by two parameters: shape parameter k and scale parameter θ . Its probability density function is given by

$$f(x; k, \theta) = \frac{1}{\theta^k \Gamma(k)} x^{k-1} e^{-\frac{x}{\theta}}$$

where $\Gamma(k) = \int_0^\infty t^{k-1} e^{-t} dt$ is the Gamma function. We now wish to find the best Gamma-distribution approximation to the true distribution of \mathcal{D} .

Your task and approach is essentially the same as in Task C. Restated for clarity, you must do the following:

- 1. First, compute an expression for the first two moments μ_1^{Gamma} , μ_2^{Gamma} of the distribution $\text{Gamma}(k, \theta)$ as a function of k and θ .
- 2. Then, use the fsolve function from scipy.optimize to compute a solution (k, θ) to $\hat{\mu}_i = \mu_i^{\text{Gamma}}$, i = 1, 2. No rounding is required, since k, θ may be real in this case. Say the found parameters are (k^*, θ^*) .
- 3. Finally, using numpy.linspace and scipy.stats.gamma.pdf (the pdf takes three parameters: find out which ones are k and θ and which one is 0), plot the binomial distribution $Bin(k^*, \theta^*)$ on top of the histogram of \mathcal{D} . You should get something like the figure in figure 4. Save the plot to 3d.png.

Solution

The mean is

$$\mu_1^{\text{Gamma}} = \int_0^\infty \frac{1}{\theta^k \Gamma(k)} x^k e^{-\frac{x}{\theta}} dx = \frac{\theta^{k+1} \Gamma(k+1)}{\theta^k \Gamma(k)} \int_0^\infty f(x; k+1, \theta) dx = k\theta \cdot 1 = k\theta,$$

since $\Gamma(k+1) = k\Gamma(k)$. Similarly, the second moment is

$$\mu_2^{\text{Gamma}} = \int_0^\infty \frac{1}{\theta^k \Gamma(k)} x^{k+1} e^{-\frac{x}{\theta}} \, \mathrm{d}x = \frac{\theta^{k+2} \Gamma(k+2)}{\theta^k \Gamma(k)} \int_0^\infty f(x;k+2,\theta) \, \mathrm{d}x = (k+1)k\theta^2.$$

The code is then similar to the first part; the only difference is that the function f now computes the error in the moments for the gamma distribution:

```
def f(x): # x = [k, theta] return [x[0]*x[1] - mean, (x[0]+1)*x[0]*x[1]**2 - second_moment]

We find (k^*, \theta^*) \approx (9.69, 0.67).

Rubric: 1.5 points each for correct first and second moments, 1 point for correct code, 1 points for the right values of (k^*, \theta^*) = (9.69, 0.67). 2 points for the correct plot.
```

 \diamond Task E (\star) [3+3]

It looks like both of the distributions did a good job, but which one did a better job?

One simple way to find out is to compute what is called the likelihood of the dataset. This is simply the probability that the dataset would be \mathcal{D} , supposing that the true distribution was actually our best-fit distribution.

Definition 4 (Likelihood). Given a dataset S and a choice of parameter $\lambda = \lambda_0$ for a family of distributions $P[\lambda]$ parameterized by λ , define the likelihood of λ_0 by

$$\mathcal{L}(\lambda_0 \mid S) := P_{\lambda_0}[S] = \prod_{i=1}^n P_{\lambda_0}[X_i].$$

Here $P_{\lambda_0}[x]$ is the PDF of the distribution whose parameter is λ_0 . In our case, $P_{\lambda}[x]$ is $Bin(\lambda = (n, p))(x)$ or $Gamma(\lambda = (k, \theta))(x)$.

For large datasets, since each probability is a small number, the likelihood can underflow. Thus, the average log-likelihood is typically calculated:

$$\ell(\theta \mid S) := \frac{\log \mathcal{L}(\theta \mid S)}{n},$$

where n is the size of the dataset. Calculate (in code, no for loops allowed) the average log-likelihood for both best-fit distributions. Since Bin(n, p)(x) is nonzero only at integer values of x, you will need to round each datapoint to the nearest integer before computing the likelihood. No such thing required for the Gamma distribution.

A larger likelihood is typically attributed to a better fit. Which distribution was a better fit?

```
Solution
```

```
# binom(n,p) case
data_rounded = np.round(data)
log_likelihood_binom = np.mean(np.log(scipy.stats.binom.pmf(data_rounded, n, p) + 1e-8))
# gamma(k,t) case
log_likelihood_gamma = np.mean(np.log(scipy.stats.gamma.pdf(data, k, 0, t)))
The binomial case yields -2.157 and the gamma case -2.161. The binomial fit is slightly better (higher
```

The binomial case yields -2.157 and the gamma case -2.161. The binomial fit is slightly better (higher likelihood).

Rubric: 3 point for correct code, 2 points for the correct log-likelihoods, 1 point for the correct conclusion.

 \diamond Task F [2+2+2+2]

Notice the two peaks in the distribution? This immediately sort of tells us that the distribution could not have been from a Binomial or Gamma function, since those have a unique mode.

A common distribution with two peaks is the Gaussian Mixture Model, which is the subject of Question 5. Please read Task A of Question 5 before moving ahead.

We are going to assume now that our distribution is composed of a two-component Gaussian mixture, each component having variance 1. That is, the distribution modelling our data is assumed to have the pdf

$$P[x] = \frac{1}{\sqrt{2\pi}} \left(p_1 \exp\left(-\frac{(x - \mu_1)^2}{2}\right) + p_2 \exp\left(-\frac{(x - \mu_2)^2}{2}\right) \right).$$

We have four parameters to find, and so need four moments. To make things easy, here are the first four moments of this distribution (here $\sigma_1 = \sigma_2 = 1$):

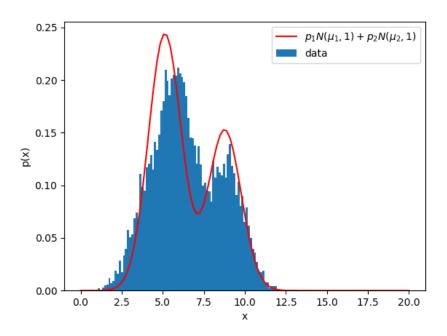


Figure 5: The best two-component unit-variance GMM approximation to the true distribution

$$\begin{split} \mu_1^{\rm gmm} &= p_1 \mu_1 + p_2 \mu_2. \\ \mu_2^{\rm gmm} &= p_1 (\sigma_1^2 + \mu_1^2) + p_2 (\sigma_2^2 + \mu_2^2). \\ \mu_3^{\rm gmm} &= p_1 (\mu_1^3 + 3\mu_1 \sigma_1^2) + p_2 (\mu_2^3 + 3\mu_2 \sigma_2^2). \\ \mu_4^{\rm gmm} &= p_1 (\mu_1^4 + 6\mu_1^2 \sigma_1^2 + 3\sigma_1^4) + p_2 (\mu_2^4 + 6\mu_2^2 \sigma_2^2 + 3\sigma_2^4). \end{split}$$

As in Tasks C and D, compute the following:

- 1. First, compute $\hat{\mu}_i$ for i = 3, 4.
- 2. Then, use the fsolve function from scipy optimize to compute a solution (μ_1, p_1, μ_2, p_2) to $\hat{\mu}_i = \mu_i^{\text{gmm}}$, i = 1, 2, 3, 4. No rounding is required. Say the found parameters are $(\mu_1^*, p_1^*, \mu_2^*, p_2^*)$.
- 3. Finally, using numpy.linspace and scipy.stats.norm.pdf, plot the GMM distribution obtained on top of the histogram of \mathcal{D} . You should get something like the figure in figure 5. Save the plot to 3f.png.

To finish, compute the average negative log-likelihood of the obtained GMM distribution. Is it better an approximation than the previous two?

```
The moments are

mu3 = np.mean(data**3)

mu4 = np.mean(data**4)

The function f is slightly more complicated, but the idea is the same.

def f(x): # x = [mu1, p1, mu2, p2]

return [

x[1]*x[0] + x[3]*x[2] - mu1,

x[1]*(1 + x[0]**2) + x[3]*(1 + x[2]**2) - mu2,

x[1]*(x[0]**3 + 3*x[0]) + x[3]*(x[2]**3 + 3*x[2]) - mu3,

x[1]*(x[0]**4 + 6*x[0]**2 + 3) + x[3]*(x[2]**4 + 6*x[2]**2 + 3) - mu4

]
```

Optimal values are $(\mu_1^*, p_1^*, \mu_2^*, p_2^*) \approx (5.13, 0.61, 8.77, 0.38)$. The log-likelihood is -2.183, which is in fact worse than the binomial and gamma approximations. However, had we added the variances as parameters as well, we would get a better fit - optimal values are $(\mu_1^*, p_1^*, \mu_2^*, p_2^*, \sigma_1^*, \sigma_2^*) \approx (5.62, 0.75, 9.18, 0.25, 1.54, 0.98)$ with a log-likelihood of -2.115, substantially better than the binomial and gamma approximations.

Rubric: 1 point for correct moment computation code, 1 point for the correct third and fourth moments, 2 points for correct further code, 1 point for the right values of $(\mu_1^*, p_1^*, \mu_2^*, p_2^*) = (5.13, 0.61, 8.77, 0.38)$, 2 points for the correct plot, 1 point for the correct log-likelihood and conclusion.

This task shows the power of a more versatile distribution like the GMM; indeed, clustering is typically done this way, using a GMM, the only difference being that an equation solver like fsolve is replaced by a heuristic called *expectation maximization*, since solvers are slow for many-component GMMs.

Teaser. A different choice of distribution, with six parameters gave the fit in figure 6. In general, an important part of machine learning is parameter estimation to fit data. Feedforward neural networks typically do just that. As we have seen, more parameters can get you closer to the "right distribution", and it is not surprising that GPT has a few billion parameters that it learns. It is not just roses, though. A large number of parameters runs the risk of overfitting. Data Analysis and ML study better ways to learn parameter values, better families of distributions, how to detect and avoid overfitting, and more.

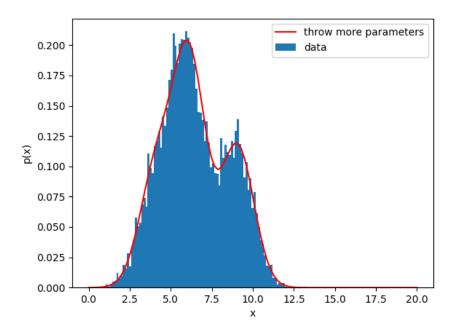


Figure 6: Two more parameters bring the total to 6 parameters estimated using the first six moments

Ungraded Bonus. Any guesses what the true distribution is? Winners can claim a coffee treat at Cafe92 from the authors.

§ 4 Quality in Inequalities

Let us dive deeper into the inequalities we have studied in class (and a new one):

Definition 5 (Markov's Inequality). *Let* X *be any non-negative random variable and* a > 0*,*

$$P[X \ge a] \le \frac{\mathbb{E}[X]}{a}.$$

♦ Task A [1+2]

Give an intuitive "proof" for this inequality and reason why it could be correct (you can try playing around with different X's and a's).

Solution

One intuitive idea is just that the mean cannot be less than the sum of a subset of members divided by the total number of members in a set (any other reasonable intuition works).

Rubric: 1 point for a reasonable intuitive proof

Now, prove this inequality rigorously for continuous random variables. Try to reason about the definition of expectation and how you can manipulate it to serve your purpose.

Solution

$$\mathbb{E}[X] = \int_0^\infty x f_X(x) dx \ge \int_a^\infty x f_X(x) dx \ge \int_a^\infty a f_X(x) dx = a P[X \ge a].$$

Rubric: 2 points for a correct proof.

♦ Task B [4]

Now that we have established this inequality, let us move on to linking it to what we have already studied in class - the Chebyshev-Cantelli inequality.

Use Markov's inequality to prove the following version of the Chebyshev-Cantelli inequality for a random variable X with mean μ and variance σ^2 : for every $\tau > 0$, we have

$$P[X - \mu \ge \tau] \le \frac{\sigma^2}{\sigma^2 + \tau^2}.$$

Solution

Let $Y = X - \mu$. $\mathbb{E}[Y] = 0$, $Var(Y) = Var(X) = \sigma^2$. Our goal becomes to show

$$P[Y \ge \tau] \le \frac{\sigma^2}{\sigma^2 + \tau^2}.$$

Using Markov's inequality directly (after squaring, of course, since Markov's inequality applies for non-negative random variables only) doesn't work (it gives Chebyshev's inequality, actually)

$$P[Y \ge \tau] \le P[Y^2 \ge \tau^2] \le \frac{\mathbb{E}[Y^2]}{\tau^2} = \frac{\sigma^2}{\tau^2} \nleq \frac{\sigma^2}{\sigma^2 + \tau^2}$$

so we need to try something else. To achieve a tighter bound, we introduce a new parameter b and consider the equivalent event $Y + b \ge \tau + b$. Again, using Markov's inequality after squaring yields

$$P[Y + b \ge \tau + b] \le P[(Y + b)^2 \ge (\tau + b)^2] \le \frac{\mathbb{E}[(Y + b)^2]}{(\tau + b)^2}.$$

the expectation of $(Y + b)^2$ is

$$\mathbb{E}\left[(Y+b)^2\right] = \mathbb{E}\left[Y^2\right] + 2b\mathbb{E}\left[Y\right] + b^2 = \mathbb{E}\left[Y^2\right] + b^2 = \sigma^2 + b^2,$$

where we used the fact that $\mathbb{E}[Y] = 0$. Thus, we have

$$P[Y \ge \tau] = P[Y + b \ge \tau + b] \le \frac{\sigma^2 + b^2}{(\tau + b)^2}.$$

As a sanity-check, note that setting b = 0 indeed yields the bound from before introducing b. Now, we need only minimize the right-hand side with respect to b to get the tightest bound on $P[Y \ge \tau]$. A derivative test yields optimal $b = \sigma^2/\tau$. Substituting this back in, we get

$$P[Y \ge \tau] \le \frac{\sigma^2 + (\sigma^2/\tau)^2}{\left(\tau + \sigma^2/\tau\right)^2} = \frac{\sigma^2}{\sigma^2 + \tau^2},$$

as desired.

Rubric: 4 points for a correct proof. Partial credit for partial progress.

 \diamond Task C(\star) [3]

Yay, our inequalities are successfully linked! Now we can move onto proving a strong bound through these inequalities...start out by showing that for a random variable X where $M_X(t)$ represents the MGF (see Question 1) for X, the following hold:

$$P[X \ge x] \le e^{-tx} M_X(t) \ \forall t > 0.$$

$$P[X \le x] \le e^{-tx} M_X(t) \ \forall t < 0.$$

Solution

The result contains the MGF $\mathbb{E}\left[e^{tX}\right]$, which makes it reasonable to use Markov's inequality on the random variable e^{tX} (it is indeed non-negative) and event $e^{tX} \ge e^{tx}$. Details follow.

For t > 0, we have $X \ge x \iff tX \ge tx \iff e^{tX} \ge e^{tx}$. Thus, we can write

$$P[X \ge x] = P\left[e^{tX} \ge e^{tx}\right] \le \frac{\mathbb{E}\left[e^{tX}\right]}{e^{tx}} = e^{-tx}M_X(t).$$

Similarly, for t < 0, we have $X \le x \iff tX \ge tx \iff e^{tX} \ge e^{tx}$. Again, we can write

$$P[X \le x] = P\left[e^{tX} \ge e^{tx}\right] \le \frac{\mathbb{E}\left[e^{tX}\right]}{e^{tx}} = e^{-tx}M_X(t).$$

Rubric: 1.5 points for each correct inequality.

 \diamond Task D (\star) [1+4+1]

Now take n **independent** Bernoulli random variables $X_1, X_2, \dots X_n$ where $\mathbb{E}[X_i] = p_i$. Since each X_i has the same distribution and is independent of all other X_j 's, we call the collection of random variables X_1, \dots, X_n a collection of *independent* and *identically distributed* (i.i.d) random variables. (they are not necessarily identically distributed in this case. However, since this line was present in the document, assuming $p_i = p$ for each i is alright.)

Let us define a new random variable Y to be the sum of these random variables, that is, $Y = \sum_{i=1}^{n} X_i$.

- 1. What is the expectation of *Y*? *Call it* μ .
- 2. Show that

$$\mathbb{P}[Y \geq (1+\delta)\mu] \leq \frac{e^{\mu(e^t-1)}}{e^{(1+\delta)t\mu}}.$$

3. Show how to improve this bound even further by choosing an appropriate value of t.

Solution

The expectation is $\mu = \mathbb{E}[Y] = \sum_{i=1}^{n} \mathbb{E}[X_i] = \sum_{i=1}^{n} p_i$. The only reasonable way to get anywhere close to an exponential on the right-hand side would be to use Task C's result to Y, which yields (for t > 0)

$$P[Y \ge (1+\delta)\mu] \le e^{-t(1+\delta)\mu} M_Y(t).$$

The magic is that $M_Y(t) = \prod_{i=1}^n M_{X_i}(t)$, which means the RHS separates neatly:

$$e^{-t(1+\delta)\mu}M_Y(t) = e^{-t(1+\delta)\mu} \prod_{i=1}^n M_{X_i}(t) = \prod_{i=1}^n e^{-t(1+\delta)p_i} M_{X_i}(t).$$

The MGF is still a bit of a mess, being

$$M_{X_i}(t) = \mathbb{E}\left[e^{tX_i}\right] = 1 + p_i(e^t - 1).$$

Here we make a clever observation: since $1 + x \le e^x$, $M_{X_i}(t) \le e^{p_i(e^t - 1)}$. Substituting this back in, we get

$$P[Y \ge (1+\delta)\mu] \le \prod_{i=1}^n e^{-t(1+\delta)p_i} e^{p_i(e^t-1)} = \prod_{i=1}^n e^{p_i(e^t-1-(1+\delta)t)} = e^{\mu(e^t-1-(1+\delta)t)}.$$

Whew, we are left with a single exponential (and the answer to the second sub-task). Finally, as with Task B, we can minimize the RHS with respect to t to get the tightest bound. Differentiating and setting to zero yields $t = \ln(1 + \delta)$, which gives the improved bound

$$P[Y \ge (1+\delta)\mu] \le \exp(\mu(\delta - (1+\delta)\ln(1+\delta))).$$

Note that we must ensure t > 0, which means the bound is only valid for $\delta > 0$. For $\delta < 0$, we can use the second part of Task C's result to get a similar bound (it is alright if omitted).

A final remark: at this value of t, the MGF approximation is very tight for small δ , since the MGF is $M_{X_i}(t) = 1 + p_i \delta \approx \exp(p_i \delta)$.

Rubric: 1 point for the expectation, 3 points for a correct derivation of the bound in sub-part 2, 1 point for the optimal value of t, 1 point for the improved bound.

The resulting best bound for $P[Y \ge (1 + \delta)\mu]$ is called a Chernoff bound and is an example of a *concentration* theorem - it can be shown that most of Y's probability density is concentrated about μ .

Chernoff bounds are related to the very useful *Central Limit Theorem* and also play very important roles in the analysis of randomized algorithms and the theory of machine learning. They are thus considered a cornerstone of probability theory. It is understood that everyone studying probability must have seen a Chernoff bound - now you know!

 \Diamond Task E [4]

Another important theorem, especially important for estimation using samples: the *weak law of large numbers* (WLLN). We shall try to prove it in this task, using the Chernoff bound.

Theorem 6. Let X_1, \dots, X_n be i.i.d random variables with each having mean μ . We define $A_n = \frac{\sum_{i=1}^n X_i}{n}$. Then for all $\epsilon > 0$, we have

$$\lim_{n\to\infty} P[|A_n - \mu| > \epsilon] = 0.$$

Essentially, the average of the variables has to be roughly constant at the value μ - it takes on any other value with probability approaching 0. Intuitively, what it means is that if you keep sampling from the same distributions and add up all of them, deviations left of the mean are cancelled by deviations right of the mean, and the net result is that the sum is always roughly the same - $n\mu$ - from which it follows that the mean is always roughly μ .

Prove WLLN using the Chernoff bound from Task D. If you did not solve Task D, you may provide a proof using just the Chebyshev inequality. However, if you did solve it, you should use the Chernoff bound obtained from Task D to prove WLLN.

Solution

Firstly, we re-write the LHS:

$$P[|A_n - \mu| > \epsilon] = P[|Y - n\mu| > n\epsilon] \le P[Y \ge \mu_Y(1 + \delta)] + P[Y \le \mu_Y(1 - \delta)],$$

where $\delta = \epsilon/\mu$ and $\mu_Y = \mathbb{E}[Y] = n\mu$. We already have a bound for the first term from Part D. The main ingredient in this proof is handling the other term similarly to Task D and then showing that both terms go to zero as $n \to \infty$. First, we show that the first term goes to zero. From Task D,

$$P[Y \ge \mu_Y(1+\delta)] \le \exp(\mu_Y(\delta - (1+\delta)\ln(1+\delta))) = \exp(n\mu(\delta - (1+\delta)\ln(1+\delta))).$$

It can be shown that for every $\delta > 0$, $\delta - (1 + \delta) \ln(1 + \delta) < 0$; in particular, at $\delta = \epsilon/\mu > 0$, its value is -k < 0. Thus,

$$0 \le P[Y \ge \mu_Y(1+\delta)] \le \exp(-kn\mu) \to 0 \text{ as } n \to \infty$$

which by Sandwich theorem implies that $P[Y \ge \mu_Y(1+\delta)] \to 0$ as $n \to \infty$. To the second term.

Notice that we cannot simply use Task D for the second term with $\delta = -\epsilon/\mu$, since that bound is only valid for $\delta > 0$. However, we can use the second part of Task C's result and the proof template of Task D to get a similar bound. First, from Task C,

$$P[Y \le \mu_Y(1-\delta)] \le e^{-t\mu_Y(1-\delta)} M_Y(t).$$

The MGF of *Y* is $M_Y(t) = \prod_{i=1}^n M_{X_i}(t) = \prod_{i=1}^n (1 + p_i(e^t - 1))$. We can use the trick as in Task D to get the same bound on the MGF:

$$M_Y(t) \le \prod_{i=1}^n e^{p_i(e^t-1)} = e^{\mu_Y(e^t-1)}.$$

Putting it together, we get the analogous Chernoff bound

$$P[Y \le \mu_Y(1-\delta)] \le e^{-t\mu_Y(1-\delta)}e^{\mu_Y(e^t-1)} = e^{\mu_Y(e^t-1-(1-\delta)t)}$$

Again, minimizing the RHS with respect to t yields $t = \ln(1 - \delta) > 0$, which gives the improved bound

$$P[Y \le \mu_Y(1-\delta)] \le \exp(\mu_Y(-\delta - (1-\delta)\ln(1-\delta))).$$

Notice this is very similar to the bound for the first term, except $\delta \to -\delta$ and the event is an upper bound on Y instead of a lower bound.

Next, it can be shown that $0 < \delta < 1$, $-\delta - (1 - \delta) \ln(1 - \delta) < 0$; in particular, when $\epsilon < \mu$, the second term goes to zero by the same argument as for the first term. For $\epsilon \ge \mu$ (so $\delta \ge 1$), we use

$$P[Y \le \mu_Y(1-\delta)] \le P[Y \le \mu_Y(1-1/2)],$$

since $\mu_Y(1-\delta) \le \mu_Y(1-1/2)$. Since $\delta = 1/2 < 1$, this upper bound is at most $\exp(-kn\mu)$ for some k > 0, which goes to zero as well.

Thus, both terms go to zero as $n \to \infty$, which completes the proof.

Proof using Chebyshev. Chernoff is not the best tool for the above; it is far too strong. The WLLN just requires an asymptotic bound, so any bound that goes to zero as $n \to \infty$ will do (it doesn't need to go *exponentially fast* to 0, however, Chernoff guarantees this as well). Chebyshev's inequality works. We have

$$P[|A_n - \mu| > \epsilon] = P[|Y - n\mu| > n\epsilon] \le \frac{\operatorname{Var}(Y)}{n^2 \epsilon^2} = \frac{n\sigma^2}{n^2 \epsilon^2} = \frac{\sigma^2}{n\epsilon^2} \to 0 \text{ as } n \to \infty.$$

It is natural to wonder why a needlessly complicated tool like Chernoff was used in the first place. This was just a trivia application to make use of the Chernoff bound; applications that actually require bounds as tight as Chernoff's are ubiquitous, however are out of scope and hence not discussed here.

Rubric: 4 points for a correct proof using Task D. If the same result in Task D was used to bound both terms, at most 1 point will be awarded. An argument of 'similar proof to Task D' for the second term will be awarded 3 points for a complete proof. Partial credit may be awarded for partial progress.

If Task D was solved, but not used, award at most 1 point (we are sorry, but you were specifically asked to use Task D. The goal of the problem was the upper bound proof of Chernoff, not really the WLLN). If Task D was not solved, a correct proof using Chebyshev fetches 2 points.

§ 5 A Pretty "Normal" Mixture

We have been looking at Gaussian (normal) random variables and their manipulation. Now we shall take many such Gaussians and *mix* them!

Definition 7 (GMM). A Gaussian Mixture Model (GMM) is a random variable defined in terms of K Gaussian random variables and follows the PDF

$$P[X = x] = \sum_{i=1}^{K} p_i P[X_i = x],$$

where each $X_i \sim \mathcal{N}(\mu_i, \sigma_i^2)$ is a Gaussian random variable with mean μ_i and variance σ_i^2 for all $i \in \{1, 2, ..., K\}$. Moreover, each $p_i \geq 0$ and $\sum_{i=1}^{K} p_i = 1$.

♦ Task A [2]

To sample from a GMM's distribution, we use the following algorithm:

- 1. First, one of the Gaussian variables X_i is randomly chosen (or effectively, an index i is chosen) according to the PMF $\{p_1, p_2, \dots p_k\}$. That is, i or X_i is chosen in this step with probability p_i .
- 2. Next, we sample a value from the chosen Gaussian's distribution $\mathcal{N}(\mu_i, \sigma_i^2)$ and this is the final value sampled from the GMM.

Suppose the output of the algorithm is random variable \mathcal{A} with PDF $f_{\mathcal{A}}$, and the PDF of the GMM variable X is f_X . Show that for every $u \in \mathbb{R}$, $f_{\mathcal{A}}(u) = f_X(u)$, that is, indeed, this algorithm samples from the GMM variable's distribution.

Solution

Let E_i denote the event that the *i*-th Gaussian is chosen. Then, we have

$$f_{\mathcal{A}}(u) = P_{\mathcal{A}}[X = u] = \sum_{i=1}^{K} P_{\mathcal{A}}[X = u | E_i] P[E_i] = \sum_{i=1}^{K} p_i P[X_i = u] = f_X(u).$$

Rubric: 2 points for a correct proof.

 \diamond Task B [1+2+2]

Let *X* be a GMM sampled by the method described above where each Gaussian $X_i \sim \mathcal{N}(\mu_i, \sigma_i^2)$ is chosen with a probability $p_i \geq 0$. Then compute

- 1. $\mathbb{E}[X]$.
- 2. Var [X].
- 3. the MGF $M_X(t)$ of X.

Solution

$$\mathbb{E}[X] = \int_{-\infty}^{\infty} x f_X(x) dx = \sum_{i=1}^{K} p_i \int_{-\infty}^{\infty} x f_{X_i}(x) dx = \sum_{i=1}^{K} p_i \mu_i.$$

$$\text{Var}(X) = \int_{-\infty}^{\infty} x^2 f_X(x) dx - \mathbb{E}[X]^2 = \sum_{i=1}^{K} p_i \int_{-\infty}^{\infty} x^2 f_{X_i}(x) dx - \mu^2 = \sum_{i=1}^{K} p_i (\sigma_i^2 + \mu_i^2) - \mu^2.$$

$$M_X(t) = \mathbb{E}\left[e^{tX}\right] = \int_{-\infty}^{\infty} e^{tx} f_X(x) dx = \sum_{i=1}^{K} p_i \int_{-\infty}^{\infty} e^{tx} f_{X_i}(x) dx = \sum_{i=1}^{K} p_i \exp\left(\mu_i t + \frac{\sigma_i^2 t^2}{2}\right).$$

Rubric: 1 point for the expectation, 2 points for the variance, 2 points for the MGF.

 \diamond Task C [1+1+2+2+1+1]

Now, we may be inclined to think "Isn't this just a weighted sum of Gaussians?" Let us now prove (or disprove) this property. Let us take a random variable *Z* to be a weighted sum of *k* **independent** Gaussian random variables,

$$Z = \sum_{i=1}^K p_i X_i,$$

where $X_i \sim \mathcal{N}(\mu_i, \sigma_i^2)$. For our new random variable Z, find the same expressions as in Task B:

- 1. $\mathbb{E}[Z]$.
- 2. Var [Z].
- 3. The PDF $f_Z(u)$ of Z.
- 4. The MGF $M_Z(t)$ of Z.
- 5. What can you conclude? Do *X* and *Z* have the same properties?
- 6. What distribution does Z seem to follow?

Solution

$$\mathbb{E}[Z] = \sum_{i=1}^K p_i \mathbb{E}[X_i] = \sum_{i=1}^K p_i \mu_i.$$

$$Var(Z) = \sum_{i=1}^{K} p_i^2 Var(X_i) = \sum_{i=1}^{K} p_i^2 \sigma_i^2,$$

where we used the independence of the X_i 's to move the variance operator inside the sum. We shall compute the MGF before the PDF.

$$M_Z(t) = \mathbb{E}\left[e^{tZ}\right] = \mathbb{E}\left[e^{t\sum_{i=1}^K p_i X_i}\right] = \prod_{i=1}^K \mathbb{E}\left[e^{tp_i X_i}\right] = \prod_{i=1}^K \exp\left(t\mu_i p_i + t^2 \frac{\sigma_i^2 p_i^2}{2}\right) = \exp\left(t\mathbb{E}\left[Z\right] + t^2 \frac{\operatorname{Var}(Z)}{2}\right).$$

The MGF is identical to that of a Gaussian random variable with mean $\mathbb{E}[Z]$ and variance Var(Z). Thus, Z is a Gaussian random variable, with PDF

$$f_Z(u) = \frac{1}{\sqrt{2\pi \text{Var}(Z)}} \exp\left(-\frac{1}{2} \left(\frac{u - \mathbb{E}[Z]}{\sqrt{\text{Var}(Z)}}\right)^2\right).$$

Clearly, $Var(X) \neq Var(Z)$, so X and Z don't have the same properties. Z follows a Gaussian distribution, as already remarked.

Rubric: 1 point for the expectation, 1 point for the variance, 2 points for the MGF, 2 points for the PDF, 1 point for the conclusion, 1 point for the distribution.

♦ Task D (B) [3]

Theorem 8. For a random variable X, if it is

- 1. either finite and discrete,
- 2. or if it is continuous and its MGF $\phi_X(t)$ is known for some (non-infinitesimal) interval,

then its MGF and PDF uniquely determine each other.

Prove the above theorem for the finite discrete case.

What can you now conclude about *X* and *Z*? Also explain logically why this may be the case.

Solution

Suppose that *X* is a finite discrete random variable taking values $\{x_1, \dots, x_n\}$. Let $M_X(t)$ be the MGF of *X* and $f_X(x)$ be the PMF of *X*. Given $f_X(x)$, $M_X(t)$ is uniquely determined by

$$M_X(t) = \mathbb{E}\left[e^{tX}\right] = \sum_{i=1}^n e^{tx_i} f_X(x_i),$$

so two random variables with the same PMF must have the same MGF. Conversely, given $M_X(t)$, we can recover $f_X(x)$ by

$$f_X(x_i) = [t^0]e^{-tx_i}M_X(t),$$

where $[t^0]g(t)$ denotes the coefficient of t^0 (constant term) in the Taylor expansion of g(t). Alternately, one can "get rid" of the non-constant exponential terms $\exp(t(x_j - x_i))$ for $j \neq i$ by going over to the complex numbers; we leave it to the reader to verify that

$$f_X(x_i) = \frac{1}{2\pi} \int_{-\pi}^{\pi} e^{-\iota t x_i} M_X(\iota t) dt.$$

Yet another way is as follows. Suppose for the sake of contradiction that there are two PMFs α and β with the same MGF. That is, for every $t \in \mathbb{R}$,

$$\sum_{i=1}^{n} e^{tx_i} \alpha(x_i) = \sum_{i=1}^{n} e^{tx_i} \beta(x_i) \iff \sum_{i=1}^{n} e^{tx_i} [\alpha(x_i) - \beta(x_i)] = 0.$$

Let $\Delta = \{i \in [n] \mid \alpha(x_i) \neq \beta(x_i)\}$; it is non-empty. Pick $i \in \Delta$ with the largest $|x_i|$. Then, for $t \to \infty$ (if $x_i > 0$) or $t \to -\infty$ (if $x_i < 0$), the exponential term e^{tx_i} dominates the sum, and so the difference between the two sides grows without bound (so it can't be zero for *every t*), a contradiction.

Finally, since $M_X(t) \neq M_Z(t)$ in general, X and Z must have different PDFs using part (2) of the theorem, and so they are different random variables (as we already concluded from Task C from their unequal variances). Another way to see this is that X has multiple modes (the μ_i 's) while Z is unimodal (since it is Gaussian), so they cannot possibly be the same random variable.

Rubric: 2 points for a correct proof. 1 point for the conclusion.