CM146, Winter 2019 Problem Set 0: Math prerequisites Due Jan 14, 2019 at 11:59 pm

Necessary Minimum Background Test [45 pts]

While you are welcome to use online resources, such as Wolfram-Alpha, you should be able to solve these problems by hand.

1 Multivariate Calculus [2 pts]

Consider $y = z \sin(x)e^{-x}$. What is the partial derivative of y with respect to x? Solution:

$$\frac{\partial y}{\partial x} = z \left(\cos(x) e^{-x} - \sin(x) e^{-x} \right)$$

2 Linear Algebra [8 pts]

Consider the matrix X and the vectors y and z below:

$$oldsymbol{X} = egin{pmatrix} 2 & 4 \\ 1 & 2 \end{pmatrix} \qquad oldsymbol{y} = egin{pmatrix} 1 \\ 3 \end{pmatrix} \qquad oldsymbol{z} = egin{pmatrix} 2 \\ 3 \end{pmatrix}$$

(a) What is the inner product $y^T z$? Solution:

$$\boldsymbol{y}^T \boldsymbol{z} = 1 \cdot 2 + 3 \cdot 3 = 11$$

(b) What is the product Xy? Solution:

$$\mathbf{X}\mathbf{y} = \begin{pmatrix} 2 \cdot 1 + 4 \cdot 3 \\ 1 \cdot 2 + 2 \cdot 3 \end{pmatrix} = \begin{pmatrix} 14 \\ 8 \end{pmatrix}$$

- (c) Is X invertible? If so, give the inverse; if not, explain why not. Solution: No. X is not invertible. The columns of X are not linearly independent: $2x_1 x_2 = 0$.
- (d) What is the rank of X? Solution: The rows (or equivalently, columns) of X are linearly dependent so rank(X) = 1.

This assignment is adapted from course material by William Cohen, Ziv Bar-Joseph (CMU) and Jessica Wu (Harvey Mudd).

Some solutions are copied verbatim or adapted from Wikipedia.

3 Probability and Statistics [10 pts]

Consider a sample of data S obtained by flipping a coin five times. $X_i, i \in \{1, ..., 5\}$ is a random variable that takes a value 0 when the outcome of coin flip i turned up heads, and 1 when it turned up tails. Assume that the outcome of each of the flips does not depend on the outcomes of any of the other flips. The sample obtained $S = (X_1, X_2, X_3, X_4, X_5) = (1, 1, 0, 1, 0)$.

- (a) What is the sample mean for this data? Solution: $\mu = \frac{1}{n} \sum_{i=1}^{n} X_i = \frac{3}{5}$
- (b) What is the unbiased sample variance? Solution: $\sigma^2 = \frac{1}{n-1} \sum_{i=1}^n (X_i \mu)^2 = \frac{1}{4} \left(3 \left(1 \frac{3}{5} \right)^2 + 2 \left(0 \frac{3}{5} \right)^2 \right)$ $\frac{1}{4} \left(3 \cdot \frac{4}{25} + 2 \cdot \frac{9}{25} \right) = \frac{3}{10}$ Give partial credit if student uses the biased estimator $-\sigma^2 = \frac{1}{n} \sum_{i=1}^n (X_i - \mu)^2$
- (c) What is the probability of observing this data assuming that a coin with an equal probability of heads and tails was used? (i.e., The probability distribution of X_i is $P(X_i = 1) = 0.5$, $P(X_i = 0) = 0.5$.) Solution: $P(S) = \left(\frac{1}{2}\right)^5 = \frac{1}{32}$
- (d) Note the probability of this data sample would be greater if the value of the probability of heads $P(X_i=1)$ was not 0.5 but some other value. What is the value that maximizes the probability of the sample S? [Optional: Can you prove your answer is correct?] **Solution:** To maximize the probability of sample S, p should take on the sample mean: $p^* = \frac{3}{5}$. Optional: Let p = P(X=1). We wish to find the value of p, $0 \le p \le 1$ that maximizes the likelihood $L(p) = p^3(1-p)^2$. To find the critical points, find p^* such that $\frac{dL}{dp}\Big|_{p^*} = 0$.

$$\frac{dL}{dp} = -2p^3(1-p) + 3p^2(1-p)^2 = p^2(1-p)(-2p+3(1-p)) = p^2(1-p)(3-5p)$$
$$p^* = 0, \frac{3}{5}, 1$$

Thus, L attains a minimum of 0 at p=0,1, and L attains a maximum of $\left(\frac{3}{5}\right)^3 \left(\frac{2}{5}\right)^2 = \frac{108}{3125}$ at $p=\frac{3}{5}$. (We could also have computed $\frac{d^2L}{dp^2}$ to determine the minimum and maximum.)

(e) Given the following joint distribution between X and Y, what is P(X = T|Y = b)?

$$\begin{array}{c|ccccc}
P(X,Y) & & Y & \\
\hline
a & b & c & \\
\hline
X & T & 0.2 & 0.1 & 0.2 & \\
F & 0.05 & 0.15 & 0.3 & \\
\end{array}$$

Solution:
$$P(X = T|Y = b) = \frac{P(X = T, Y = b)}{P(Y = b)} = \frac{0.1}{0.1 + 0.15} = 0.4.$$

4 Probability axioms [5 pts]

Let A and B be two discrete random variables. In general, are the following true or false? (Here A^c denotes complement of the event A.)

(a)
$$P(A \cup B) = P(A \cap (B \cap A^c))$$
 Solution: False. $P(A \cap (B \cap A^c)) = P((A \cap A^c) \cap B) = 0$.

- (b) $P(A \cup B) = P(A) + P(B)$ Solution: True. This is the inclusion-exclusion rule.
- (c) $P(A) = P(A \cap B) + P(A^c \cap B)$ Solution: False. $P(A \cap B) + P(A^c \cap B) = P(B)$.
- (d) P(A|B) = P(B|A) Solution: False. Conditional probabilities are not commutative. [Compare $P(A|B) = \frac{P(A \cap B)}{P(B)}$ against $P(B|A) = \frac{P(A \cap B)}{P(A)}$.]
- (e) $P(A_1 \cap A_2 \cap A_3) = P(A_3 | (A_2 \cap A_1)) P(A_2 | A_1) P(A_1)$ Solution: True. $P(A_3 | (A_2 \cap A_1)) P(A_2 | A_1) P(A_1) = \frac{P(A_1 \cap A_2 \cap A_3)}{P(A_1 \cap A_2)} \cdot \frac{P(A_1 \cap A_2)}{P(A_1)} \cdot P(A_1) = P(A_1 \cap A_2 \cap A_3).$

5 Discrete and Continuous Distributions [5 pts]

Match the distribution name to its formula.

- (i) $p^x(1-p)^{1-x}$, when $x \in \{0,1\}; 0$ otherwise (a) Gaussian
- (ii) $\frac{1}{b-a}$ when $a \le x \le b$; 0 otherwise (iii) $\binom{n}{x} p^x (1-p)^{n-x}$ (b) Exponential
- (c) Uniform
- (d) Bernoulli
- (iv) $\lambda e^{-\lambda x}$ when $x \ge 0$; 0 otherwise (v) $\frac{1}{\sqrt{(2\pi)\sigma^2}} \exp\left(-\frac{1}{2\sigma^2}(x-\mu)^2\right)$ (e) Binomial

Solution: a-v, b-iv, c-ii, d-i, e-iii

Mean and Variance[5 pts] 6

(a) What is the mean and variance of a Bernoulli(p) random variable? Solution: The PMF for the Bernouilli distribution is

$$f(k; p) = \begin{cases} p & \text{if } k = 1\\ 1 - p & \text{if } k = 0 \end{cases}$$
$$= p^k (1 - p)^{1 - k} \text{ for } k \in \{0, 1\}$$

Let q = 1 - p. Then

$$\mu = \mathbb{E}[X] = \sum_{k \in \{0,1\}} kf(k;p) = (0)(q) + (1)(p) = p$$

$$\sigma^2 = \operatorname{Var}(X) = \mathbb{E}\left[(X - \mu)^2\right] = \mathbb{E}\left[X^2\right] - (\mathbb{E}[X])^2 = (0)^2(q) + (1)^2(p) - (p)^2 = p(1 - p) = pq$$

$$H = \mathbb{E}\left[-\ln(f(k;p))\right] = (-\ln q)(q) + (-\ln p)(p) = -q\ln q - p\ln p$$

(b) If the variance of a zero-mean random variable X is σ^2 , what is the variance of 2X? What about the variance of X + 3? Solution: If all values are scaled by a constant, the variance is scaled by the square of that constant: $Var(aX) = a^2Var(X)$. Thus, $Var(2X) = 4Var(X) = a^2Var(X)$ $4\sigma^2$.

Variance is invariant with respect to changes in a location parameter; that is, if a constant is added to all values of the variable, the variance is unchanged: Var(X + a) = Var(X). Thus, $Var(X+3) = Var(X) = \sigma^2$.

7 Algorithms [10 pts]

(a) Big-O notation

For each pair (f,g) of functions below, list which of the following are true: f(n) = O(g(n)), g(n) = O(f(n)), or both. Briefly justify your answers.

- i. f(n) = ln(n), g(n) = lg(n). Note that ln denotes log to the base e and lg denotes log to the base 2. **Solution:** Both since both functions are equivalent upto a multiplicative constant
- ii. $f(n) = 3^n, g(n) = n^{10}$ Solution: g(n) = O(f(n)) since f(n) grows much more rapidly as n becomes large.
- iii. $f(n) = 3^n, g(n) = 2^n$ Solution: g(n) = O(f(n)) since f(n) grows much more rapidly as n becomes large.

(b) Divide and Conquer

Assume that you are given an array with n elements all entries equal either to 0 or +1 such that all 0 entries appear before +1 entries. You need to find the index where the transition happens, *i.e.*, you need to report the index with the last occurrence of 0. Give an algorithm that runs in time $O(\log n)$. Explain your algorithm in words, describe why the algorithm is correct, and justify its running time. **Solution:** We give an algorithm below, called find-transition, for this problem:

find-transition (i , j) let mid = i + floor ((j-i) /2) let a = array [mid] and b = array [mid + 1] if (a == b == 1) return find-transition (i , mid) else if (a == b == 0) return find-transition (mid + 1 , j) else /* a == 0 and b == 1 */ return mid (last occurrence of 0)

This algorithm returns the correct results: note that for each recursive call we know that array[i]==0 and array[j]==1. When we stop, we know that a==0 and b==1, so we array[mid] is the last entry with 0. The running time can be analyzed via the recurrence: T(n) = T(n/2) + O(1), T(n) = c for n-4 which solves to $O(\log n)$. This algorithm is based on the idea of binary search and hence the running time is as expected.

Moderate Background Test [35 pts]

8 Probability and Random Variables [5 pts]

(a) Mutual and Conditional Independence If X and Y are independent random variables, show that $\mathbb{E}[XY] = \mathbb{E}[X]E[Y]$. Solution: If X and Y are independent, then by definition j(x,y) = f(x)g(y) where j is the joint PDF for X and Y and g are the marginal PDFs for X and Y, respectively. Then

$$\mathbb{E}[XY] = \iint xy \, j(x,y) \, dx \, dy = \iint xy \, f(x)g(y) \, dx \, dy = \left[\int xf(x) \, dx \right] \left[\int yf(y) \, dy \right] = \mathbb{E}[X]E[Y]$$

(b) Law of Large Numbers and Central Limit Theorem

Provide one line justifications.

- i. If a fair die is rolled 6000 times, the number of times 3 shows up is close to 1000. **Solution:** Assuming a fair die, the number of times 3 shows up should be close to 1000 due to the Law of Large Numbers.
- ii. If a fair coin is tossed n times and \bar{X} denotes the average number of heads, then the distribution of \bar{X} satisfies

$$\sqrt{n}(\bar{X} - \frac{1}{2}) \xrightarrow{n \to \infty} \mathcal{N}(0, \frac{1}{4})$$

Solution: The expression on the left-hand-side should tend to the expression on the right-hand-side as $n \to \infty$ due to the Central Limit Theorem.

9 Linear Algebra [20 pts]

(a) Vector Norms [4 pts]

Draw the regions corresponding to vectors $x \in \mathbb{R}^2$ with following norms (you can hand draw or use software for this question): Solution: Each l_p -norm corresponds to a unit ball. In \mathbb{R}^2 , this corresponds to a disc of diameter 2 centered at the origin.

- i. $||x||_2 \le 1$ (Recall $||x||_2 = \sqrt{\sum_i x_i^2}$.) Solution: the unit circle (origin-centered circle with diameter of 2), and its interior
- ii. $||x||_0 \le 1$ (Recall $||x||_0 = \sum_{i:x_i \ne 0} 1$.) Solution: the axes
- iii. $||\boldsymbol{x}||_1 \leq 1$ (Recall $||\boldsymbol{x}||_1 = \sum_i |x_i|$.) Solution: origin-centered diamond with diameter of 2, and its interior
- iv. $||x||_{\infty} \leq 1$ (Recall $||x||_{\infty} = \max_i |x_i|$.) Solution: origin-centered square with side length of 2, and its interior

(b) Matrix Decompositions [6 pts]

i. Give the definition of the eigenvalues and the eigenvectors of a square matrix. Solution: An eigenvector of a square matrix is a vector that points in a direction that is invariant under the associated linear transformation. In other words, if v is a vector that is not

zero, then it is an eigenvector of a square matrix A if Av is a scalar multiple of v. This condition could be written as the equation

$$Av = \lambda v$$

where λ is a number (scalar) known as the *eigenvalue* associated with the eigenvector \boldsymbol{v} .

ii. Find the eigenvalues and eigenvectors of

$$A = \begin{pmatrix} 2 & 1 \\ 1 & 2 \end{pmatrix}$$
.

Solution: The eigenvectors \boldsymbol{v} of this transformation satisfy the equation $\boldsymbol{A}\boldsymbol{v} = \lambda \boldsymbol{v}$. Rearrange this equation to obtain $(\boldsymbol{A} - \lambda \boldsymbol{I})\boldsymbol{v} = 0$, which has a solution only when its determinant $|\boldsymbol{A} - \lambda \boldsymbol{I}|$ equals zero. Set the determinant to zero to obtain the polynomial equation,

$$p(\lambda) = |\mathbf{A} - \lambda \mathbf{I}| = (2 - \lambda)^2 - 1 = (\lambda^2 - 4\lambda + 4) - 1 = \lambda^2 - 4\lambda + 3 = (\lambda - 3)(\lambda - 1) = 0,$$

known as the characteristic polynomial of the matrix A. In this case, it has the roots $\lambda = 1$ and $\lambda = 3$.

For $\lambda = 1$, the equation becomes,

$$(\boldsymbol{A} - \boldsymbol{I})\boldsymbol{v} = \begin{pmatrix} 1 & 1 \\ 1 & 1 \end{pmatrix} \begin{pmatrix} v_1 \\ v_2 \end{pmatrix} = \begin{pmatrix} 0 \\ 0 \end{pmatrix},$$

which has the solution,

$$v = \begin{pmatrix} 1 \\ -1 \end{pmatrix}$$
.

For $\lambda = 3$, the equation becomes,

$$(\mathbf{A} - 3\mathbf{I})\mathbf{w} = \begin{pmatrix} -1 & 1 \\ 1 & -1 \end{pmatrix} \begin{pmatrix} w_1 \\ w_2 \end{pmatrix} = \begin{pmatrix} 0 \\ 0 \end{pmatrix},$$

which has the solution,

$$\boldsymbol{w} = \begin{pmatrix} 1 \\ 1 \end{pmatrix}$$
.

Thus, the vectors \boldsymbol{v} and \boldsymbol{w} are eigenvectors of \boldsymbol{A} associated with the eigenvalues $\lambda=1$ and $\lambda=3$, respectively.

iii. For any positive integer k, show that the eigenvalues of \mathbf{A}^k are $\lambda_1^k, \lambda_2^k, \ldots, \lambda_n^k$, the k^{th} powers of the eigenvalues of matrix \mathbf{A} , and that each eigenvector of \mathbf{A} is still an eigenvector of \mathbf{A}^k . Solution: We will prove by induction.

Base case (given): $Av = \lambda v$

Inductive step: Prove that if $\mathbf{A}^k \mathbf{v} = \lambda^k \mathbf{v}$, then $\mathbf{A}^{k+1} \mathbf{v} = \lambda^{k+1} \mathbf{v}$.

This holds because $\mathbf{A}^{k+1}\mathbf{v} = \mathbf{A}(\mathbf{A}^k\mathbf{v}) = \mathbf{A}(\lambda^k\mathbf{v}) = \lambda^k(\mathbf{A}\mathbf{v}) = \lambda^{k+1}\mathbf{v}$.

Since both the basis and the inductive step have been performed, by mathematical induction, the statement $\mathbf{A}^k \mathbf{v} = \lambda^k \mathbf{v}$ holds for all positive integer k. Q.E.D.

(c) Vector and Matrix Calculus [5 pts]

Consider the vectors x and a and the symmetric matrix A.

- i. What is the first derivative of $\boldsymbol{a}^T\boldsymbol{x}$ with respect to \boldsymbol{x} ? Solution: Let a_i and x_i denote the elements of \boldsymbol{a} and \boldsymbol{x} , respectively. Then $f(\boldsymbol{x}) = \boldsymbol{a}^T\boldsymbol{x} = \sum_{i=1}^n a_i x_i$ so $\frac{\partial f(\boldsymbol{x})}{\partial x_k} = \frac{\partial}{\partial x_k} \sum_{i=1}^n a_i x_i = a_k$ and $\nabla_{\boldsymbol{x}} \boldsymbol{a}^T \boldsymbol{x} = \boldsymbol{a}$.
- ii. What is the first derivative of $\mathbf{x}^T \mathbf{A} \mathbf{x}$ with respect to \mathbf{x} ? What is the second derivative? Solution: Let a_{ij} denote the element in row i, column j of \mathbf{A} . Then $f(\mathbf{x}) = \mathbf{x}^T \mathbf{A} \mathbf{x} = \sum_{i=1}^n \sum_{j=1}^n x_i a_{ij} x_j$ so

$$\begin{split} \frac{\partial f(\boldsymbol{x})}{\partial x_k} &= \frac{\partial}{\partial x_k} \sum_{i=1}^n \sum_{j=1}^n a_{ij} x_i x_j \\ &= \frac{\partial}{\partial x_k} \left[\sum_{i \neq k} \sum_{j \neq k} a_{ij} x_i x_j + \sum_{i \neq k} a_{ik} x_i x_k + \sum_{j \neq k} a_{kj} x_k x_j + a_{kk} x_k^2 \right] \\ &= 0 + \sum_{i \neq k} a_{ik} x_i + \sum_{j \neq k} a_{kj} x_j + 2a_{kk} x_k \\ &= \sum_{i=1}^n a_{ik} x_i + \sum_{j=1}^n a_{kj} x_j = 2 \sum_{i=1}^n a_{ki} x_i \end{split}$$

where the last equality follows from $a_{ij} = a_{ji}$. Note that the k^{th} entry of $\nabla_{\boldsymbol{x}} f(\boldsymbol{x})$ is just the inner product of the k^{th} row of \boldsymbol{A} and \boldsymbol{x} ; therefore, $\nabla_{\boldsymbol{x}} \boldsymbol{x}^T \boldsymbol{A} \boldsymbol{x} = 2 \boldsymbol{A} \boldsymbol{x}$.

To find the second derivative,

$$\frac{\partial^2 f(\mathbf{x})}{\partial x_k x_\ell} = \frac{\partial}{\partial x_k} \left[\frac{\partial f(\mathbf{x})}{\partial x_\ell} \right] = \frac{\partial}{\partial x_k} \left[2 \sum_{i=1}^n a_{\ell i} x_i \right] = 2a_{\ell k} = 2a_{k\ell}$$

Therefore, $\nabla_{\boldsymbol{x}}^2 \boldsymbol{x}^T \boldsymbol{A} \boldsymbol{x} = 2 \boldsymbol{A}$.

(d) Geometry [5 pts]

i. Show that the vector \boldsymbol{w} is orthogonal to the line $\boldsymbol{w}^T\boldsymbol{x}+b=0$. (Hint: Consider two points $\boldsymbol{x}_1,\boldsymbol{x}_2$ that lie on the line. What is the inner product $\boldsymbol{w}^T(\boldsymbol{x}_1-\boldsymbol{x}_2)$?) Solution: This line is all \boldsymbol{x} such that $\boldsymbol{w}^T\boldsymbol{x}+b=0$. Consider two such points, called \boldsymbol{x}_1 and \boldsymbol{x}_2 . Note that $\boldsymbol{x}_1-\boldsymbol{x}_2$ is a vector parallel to our line. Also note that

$$\boldsymbol{w}^T \boldsymbol{x}_1 + b = 0 = \boldsymbol{w}^T \boldsymbol{x}_2 + b \Longrightarrow \boldsymbol{w}^T \boldsymbol{x}_1 = \boldsymbol{w}^T \boldsymbol{x}_2 \Longrightarrow \boldsymbol{w}^T (\boldsymbol{x}_1 - \boldsymbol{x}_2) = 0,$$

which shows that the vector \boldsymbol{w} is orthogonal to our line.

ii. Argue that the distance from the origin to the line $\mathbf{w}^T \mathbf{x} + b = 0$ is $\frac{b}{\|\mathbf{w}\|_2}$. Solution: We can show this by first finding the closest point to the origin that lies on this line, and then finding the distance to this point. Let \mathbf{a}^* be the closest point to the origin the lies on the line. We can write \mathbf{a}^* as

$$\mathbf{a}^* = \min_{\mathbf{s.t.}} \mathbf{a}^T \mathbf{a}$$

s.t. $\mathbf{w}^T \mathbf{a} + b = 0$

So we first solve this constrained optimization problem and find a^* . We start by taking the derivative of the objective, setting it to zero, and using Lagrange multipliers (i.e. setting the derivative of the Lagrangian $a^T a - \lambda(w^T a + b)$ to zero). We can write

$$\nabla_{\boldsymbol{a}} \left[\boldsymbol{a}^T \boldsymbol{a} - \lambda (\boldsymbol{w}^T \boldsymbol{a} + b) \right] = 2\boldsymbol{a} - \lambda \boldsymbol{w} = 0$$

to find that $a^* = \frac{\lambda}{2} w$. Hence, plugging this value for a into the constraint, we can write

$$\boldsymbol{w}^T \boldsymbol{a} + b = \boldsymbol{w}^T \left(\frac{\lambda}{2} \boldsymbol{w}\right) + b = \frac{\lambda}{2} \boldsymbol{w}^T \boldsymbol{w} + b = 0$$

Thus,

$$\lambda = \frac{-2b}{\boldsymbol{w}^T \boldsymbol{w}}$$
$$\boldsymbol{a}^* = \frac{-b}{\boldsymbol{w}^T \boldsymbol{w}} \boldsymbol{w}$$

Once we have a^* , we can compute the distance between a^* and the origin to get

$$distance = ||\boldsymbol{a}|| = \sqrt{\left(\boldsymbol{a}^*\right)^T\boldsymbol{a}^*} = \sqrt{\left(\frac{-b}{\boldsymbol{w}^T\boldsymbol{w}}\right)^2\boldsymbol{w}^T\boldsymbol{w}} = \frac{b}{\boldsymbol{w}^T\boldsymbol{w}}\sqrt{\boldsymbol{w}^T\boldsymbol{w}} = \frac{b}{\sqrt{\boldsymbol{w}^T\boldsymbol{w}}} = \frac{b}{||\boldsymbol{w}||}$$

Programming Skills

Start familiarizing yourself with the Python libraries numpy and matplotlib by completing the following exercises. (You do not have to submit your code.)

You may find the following references helpful:

- http://docs.scipy.org/doc/numpy/reference/generated/numpy.random.multivariate_normal.html
- http://docs.scipy.org/doc/numpy/reference/generated/numpy.linalg.eig.html

10 Sampling from a Distribution [2.5 pts]

For questions (a-e), only submit your plots. You do not need to submit code.

- (a) Draw 1000 samples $\mathbf{x} = \begin{pmatrix} x_1 \\ x_2 \end{pmatrix}$ from a 2-dimensional Gaussian distribution with mean $\begin{pmatrix} 0 \\ 0 \end{pmatrix}$ and identity covariance matrix, i.e. $p(\mathbf{x}) = \frac{1}{\sqrt{(2\pi)^2}} \exp\left(-\frac{||\mathbf{x}||^2}{2}\right)$, and make a scatter plot $(x_1 \text{ vs } x_2)$. Solution: See attached code.
- (b) How does the scatter plot change if the mean is $\binom{1}{1}$? **Solution:** See solution code. The data moves up and to the left. Namely, the center of the data moves from roughly [0,0] to roughly [-1,1].
- (c) How does the (original) scatter plot change if you double the variance of each component? **Solution:** See solution code. The data become more "spread out".
- (d) How does the (original) scatter plot change if the covariance matrix is changed to $\binom{1}{0.5} \binom{0.5}{1}$? **Solution:** See solution code. The data become skewed so that they stretch from the lower left to the upper right.
- (e) How does the (original) scatter plot change if the covariance matrix is changed to $\begin{pmatrix} 1 & -0.5 \\ -0.5 & 1 \end{pmatrix}$? **Solution:** See solution code. The data become skewed so that they stretch from the upper left to the bottom right.

11 Eigendecomposition [2.5 pts]

Write a python program to compute the eigenvector corresponding to the largest eigenvalue of the following matrix and submit the computed eigenvector.

$$\mathbf{A} = \begin{pmatrix} 1 & 0 \\ 1 & 3 \end{pmatrix}$$

. **Solution:** See solution code. $\lambda_1=3, \boldsymbol{v}_1=\left(\begin{smallmatrix}0\\1\end{smallmatrix}\right)$. Optional: The other eigenvalue and eigenvector is $\lambda_2=1, \boldsymbol{v}_2=\left(\begin{smallmatrix}0.894\\-0.447\end{smallmatrix}\right)$.

12 Data [5 pts]

There are now lots of really interesting data sets publicly available to play with. They range in size, quality and the type of features and have resulted in many new machine learning techniques being developed.

Find a public, free, supervised (i.e. it must have features and labels), machine learning dataset. You may NOT list a data set from 1) The UCI Machine Learning Repository or 2) from Kaggle.com. Once you have found the data set, provide the following information:

- (a) The name of the data set.
- (b) Where the data can be obtained.
- (c) A brief (i.e. 1-2 sentences) description of the data set including what the features are and what is being predicted.
- (d) The number of examples in the data set.
- (e) The number of features for each example. If this is not concrete (i.e. it is text), then a short description of the features.

For this question, do not just copy and paste the description from the website or the paper; reference it, but use your own words. Your goal here is to convince the staff that you have taken the time to understand the data set, where it came from, and potential issues involved.