## Homeworks

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## Chapter 1

## Number systems

### 1.1 Problems NS-1

## Topic of this homework:

Introduction to Matlab/Octave (see the Matlab or Octave tutorial for help)
Deliverables: Report with charts and answers to questions.

## Plotting complex quantities in Octave/Matlab

Problem \# 1: Consider the functions $f(s)=s^{2}+6 s+25$ and $g(s)=s^{2}+6 s+5$.

- 1.1: Find the zeros of functions $f(s)$ and $g(s)$ using the command roots ().

Ans:

- 1.2: Show the roots of $f(s)$ as red circles and of $g(s)$ as blue plus signs.

The $x$-axis should display the real part of each root, and the $y$-axis should display the imaginary part. Use hold on and grid on when plotting the roots.

Ans:

- 1.3 Give your figure the title "Complex Roots of $f(s)$ and $g(s)$." Label the $x$ - and $y$-axes "Real Part" and "Imaginary Part." Hint: Use xlabel, ylabel, ylim([-10 10]), and xlim([-10 10]) to expand the axes.

Problem \# 2: Consider the function $h(t)=e^{2 \pi f t}$ for $f=5$ and $t=[0: 0.01: 2]$.

- 2.1: Use subplot to show the real and imaginary parts of $h(t)$.

Make two graphs in one figure. Label the $x$-axes "Time (s)" and the $y$-axes "Real Part" and "Imaginary Part."
Ans:

- 2.2: Use subplot to plot the magnitude and phase parts of $h(t)$.

Use the command angle or unwrap (angle ()) to plot the phase. Label the $x$-axes "Time (s)" and the $y$-axes ' 'Magnitude" and "Phase (radians)."

Ans:

## Prime numbers, infinity, etc. in Octave/Matlab

Problem \# 3: Prime numbers, infinity, etc.

- 3.1: Use the Matlab/Octave function factor to find the prime factors of 123, 248, 1767, and 999,999.
Ans:
- 3.2: Use the Matlab/Octave function isprime to check if 2, 3 and 4 are prime numbers. What does the function isprime return when a number is prime, or not prime? Why?

Ans:

- 3.3: Use the Matlab/Octave function primes.m to generate prime numbers between 1 and $10^{6}$ Save them in a vector x . Plot this result using the command hist (x). Ans:
- 3.4: Now try [n,bincenters] $=$ hist $(x)$.

Use length ( n ) to find the number of bins. Ans:

- 3.5: Set the number of bins to 100 by using an extra input argument to the function hist. Show the resulting figure and give it a title and axes labels. Ans:

Problem \# 4: Inf, NaN and logarithms in Octave/Matlab

- 4.1: Try $1 / 0$ and $0 / 0$ in the Octave/Matlab command window.

What are the results? What do these 'numbers' mean in Octave/Matlab? Ans:
-4.2: Try $\log (0), \log 10(0)$ and $\log 2(0)$ in the command window.
In Matlab/Octave, the natural logarithm $\ln (\cdot)$ is computed using the function $\log$. Functions $\log _{10}$, and $\log _{2}$ are computed using $\log 10$ and $\log 2$. Ans:
-4.3: Try $\log (1)$ in the command window. What you expect for $\log 10(1)$ and $\log 2(1)$ ? Ans:
-4.4: Try $\log (-1)$ in the command window. What do you expectfor $\log 10(-1)$ and $\log 2(-1)$ ? Ans:

- 4.5: Show how Matlab/Octave arrives at the above answer because $-1=e^{i \pi}$.

Ans:

- 4.6: Try $\log (\exp (j * \operatorname{sqrt}(\mathrm{pi})))$ (i.e., $\left.\log e^{\sqrt{\pi}}\right)$ in the command window. What do you expect?
Ans:
- 4.7: What does inverse mean in this context? What is the inverse of $\ln f(x)$ ?


## Ans:

- 4.8: What is a decibel? (Look up decibels on the internet.)


## Ans:

Problem \# 5: Very large primes on Intel computers

- 5.1: Find the largest prime number that can be stored on an Intel 64 bit computer, which we call $\pi_{\text {max }}$.
Hint: As explained in the Matlab/Octave command help flintmax, the largest positive integer is $2^{53}$, however the largest integer that can be factored is $2^{32}=\sqrt{2^{64}}$. Explain the logic of your answer. Hint: help isprime(). Ans:

Problem \# 6: Suppose you are interested in primes that are greater than $\pi_{\max }$. How can you find them on an Intel computer (i.e., one using IEEE-floating point)?

- 6.1: Extending the number of primes you may considered.

Hint 1: Use uint 64 (myprimes) to extend the numbers unsigned 64 bit integers (we don't need negative primes). Hint 2: Since every prime number greater than 2 is odd, there is no reason to check the even numbers. Starting from 3 (not 2 ). $n_{\text {odd }} \in \mathbb{N} / 2$ contain all the primes other than 2 . Ans:

Problem \# 7: The following idenity is interesting:

$$
\begin{aligned}
1 & =1^{2} \\
1+3 & =2^{2} \\
1+3+5 & =3^{2} \\
1+3+5+7 & =4^{2} \\
1+3+5+7+9 & =5^{2} \\
\cdots & \\
\sum_{n=0}^{N-1} 2 n+1 & =N^{2} .
\end{aligned}
$$

- 7.1: Can you find a proof? ${ }^{1}$

Ans:

[^0]
### 1.2 Problems NS-2

## Topic of this homework:

Prime numbers, greatest common divisors, the continued fraction algorithm

## Prime numbers

Problem \# 1: Every integer may be written as a product of primes.

- 1.1: Write the numbers $1,000,000,1,000,004$, and 999,999 in the form $N=\prod_{k} \pi_{k}^{\beta_{k}}$. Hint: Use Matlab/Octave to find the prime factors.
Ans:
- 1.2: Give a generalized formula for the natural logarithm of a number $\ln (N)$ in terms of its primes $\pi_{k}$ and their multiplicities $\beta_{k}$. Express your answer as a sum of terms.


## Ans:

Problem \# 2: Using the computer

- 2.1: Explain why the following brief Matlab/Octave program returns the prime numbers $\pi_{k}$ between 1 and 100.

```
n=2:100;
k = isprime(n);
n(k)
Ans:
```

- 2.2: How many primes are there between 2 and $N=100$ ?


## Ans:

Problem \# 3: Prime numbers may be identified using a sieve (See Figure).

- 3.1: By hand, complete the sieve of Eratosthenes for $n=1, \ldots, 49$. Circle each prime $p$, then cross out each number that is a multiple of $p$.

| (1) | (2) | $(3)$ | 4 | 5 | 6 | 7 | 8 | 9 | 10 |
| :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- |
| $(11)$ | 11 | $(13$ | 14 | 15 | 16 | 17 | 18 | 19 | 20 |
| 21 | 22 | $(23)$ | 24 | 25 | 26 | 27 | 28 | 29 | 30 |
| $(31)$ | 32 | 33 | 34 | 35 | 36 | 37 | 38 | 39 | 40 |
| $(41)$ | 42 | 43 | 44 | 45 | 46 | 47 | 48 | 49 |  |

Note: 1 should not be circled as it is not a prime.

- 3.2: What is the largest number you need to consider before only primes remain?


## Ans:

-3.3: Generalize: For $n=1, \ldots, N$, what is the largest number you need to consider before only the primes remain?
Ans:

- 3.4: Write each of these numbers as a product of primes: $22,30,34,43,44,48,49$. Ans:
-3.5: Find the largest prime $\pi_{k} \leq 100$. Do not use Matlab/Octave other than to check your answer. Hint: Write the numbers starting with 100 and count backward: 100, 99, 98, 97, . . . Cross off the even numbers, leaving $99,97,95, \ldots$ Pull out a factor (only one is necessary to show that it is not prime).
Ans:
- 3.6: Find the largest prime $\pi_{k} \leq 1000$. Do not use Matlab/Octave other than to check your answer.
Ans:
- 3.7: Explain why $\pi_{k}^{-s}=e^{-s \ln \pi_{k}}$.

Ans:

## Greatest common divisors

Consider using the Euclidean algorithm to find the greatest common divisor (i.e., GCD; the largest common prime factor) of two numbers. Note that this algorithm may be performed using one of two methods:

| Method | Division | Subtraction |
| :--- | :--- | :--- |
| On each iteration... | $a_{i+1}=b_{i}$ | $a_{i+1}=\max \left(a_{i}, b_{i}\right)-\min \left(a_{i}, b_{i}\right)$ |
|  | $b_{i+1}=a_{i}-b_{i} \cdot$ floor $\left(a_{i} / b_{i}\right)$ | $b_{i+1}=\min \left(a_{i}, b_{i}\right)$ |
| Terminates when... | $b=0(\mathrm{GCD}=a)$ | $b=0(\mathrm{GCD}=a)$ |

The division method (Eq. 2.1, Sec. 2.1.2, Ch. 2) is preferred because the subtraction method is much slower.
Problem \# 4: Understanding the Euclidean algorithm (GCD)
4.1: Use the Octave/Matlab command factor to find the prime factors of $a=85$ and $b=15$. Ans:
-4.2: What is the greatest common prime factor of $a=85$ and $b=15$ ?
Ans:
-4.3: By hand, perform the Euclidean algorithm for $a=85$ and $b=15$.

## Ans:

- 4.4: By hand, perform the Euclidean algorithm for $a=75$ and $b=25$. Is the result a prime number?
Ans:
-4.5: Consider the first step of the GCD division algorithm when $a<b$ (e.g., $a=25$ and $b=75$ ). What happens to $a$ and $b$ in the first step? Does it matter if you begin the algorithm with $a<b$ rather than $b<a$ ?
Ans:
-4.6: Describe in your own words how the GCD algorithm works. Try the algorithm using numbers that have already been divided into factors (e.g., $a=5 \cdot 3$ and $b=7 \cdot 3$ ).
Ans:
- 4.7: Find the GCD of $2 \cdot \pi_{25}$ and $3 \cdot \pi_{25}$.

Ans: .

Problem \# 5: Coprimes

- 5.1: Define the term coprime.

Ans:

- 5.2: How can the Euclidean algorithm be used to identify coprimes?


## Ans:

## - 5.3: Give at least one application of the Euclidean algorithm.

## Ans:

- 5.4: Write a Matlab function, function $x=\operatorname{myg} g c d(a, b)$, that uses the Euclidean algorithm to find the GCD of any two inputs a and b. Test your function on the $(a, b)$ combinations from the previous problem. Include a printout (or hand-write) your algorithm to turn in.
Hints and advice:
- Don't give your variables the same names as Matlab functions! Since gcd is an existing Matlab/Octave function, if you use it as a variable or function name, you won't be able to use gcd to check your gcd () function. Try clear all to recover from this problem.
- Try using a "while" loop for this exercise (see Matlab documentation for help).
- You may need to use some temporary variables for a and b in order to perform the algorithm.


## Ans:

## Algebraic generalization of the GCD (Euclidean) algorithm

Problem \# 6: In this problem we are looking for integer solutions $(m, n) \in \mathbb{Z}$ to the equations $m a+n b=g c d(a, b)$ and $m a+n b=0$ given positive integers $(a, b) \in \mathbb{Z}^{+}$.
Note that this requires that either $m$ or $n$ be negative. These solutions may be found using the Euclidean algorithm only if $(a, b)$ are coprime $(a \perp b)$. Note that integer (whole number) polynomial relations such as these are known as Diophantine equations. Such equations (e.g., $m a+n b=0$ ) are linear Diophantine equations, possibly the simplest form of such relations.

Example: $\operatorname{gcd}(2,3)=1$ : $\quad \operatorname{For}(a, b)=(2,3)$, the result is

$$
\left[\begin{array}{l}
1 \\
0
\end{array}\right]=\left[\begin{array}{cc}
0 & 1 \\
1 & -2
\end{array}\right]\left[\begin{array}{cc}
0 & 1 \\
1 & -1
\end{array}\right]\left[\begin{array}{ll}
0 & 1 \\
1 & 0
\end{array}\right]\left[\begin{array}{l}
2 \\
3
\end{array}\right]=\underbrace{\left[\begin{array}{cc}
-1 & 1 \\
3 & -2
\end{array}\right]}_{m}\left[\begin{array}{l}
2 \\
3
\end{array}\right] .
$$

Thus from the above equation we find the solution $(m, n)$ to the integer equation

$$
2 m+3 n=\operatorname{gcd}(2,3)=1
$$

namely, $(m, n)=(-1,1)$ (i.e., $-2+3=1$ ). There is also a second solution $(3,-2)$ (i.e., $3 \cdot 2-2 \cdot 3=0)$ that represents the terminating condition. Thus these two solutions are a pair and the solution exists only if ( $a, b$ ) are coprime ( $a \perp b$ ).

Subtraction method: This method is more complicated than the division algorithm because at each stage we must check whether $a<b$. Define

$$
\left[\begin{array}{l}
a_{0} \\
b_{0}
\end{array}\right]=\left[\begin{array}{l}
a \\
b
\end{array}\right], \quad Q=\left[\begin{array}{cc}
1 & -1 \\
0 & 1
\end{array}\right], \quad S=\left[\begin{array}{ll}
0 & 1 \\
1 & 0
\end{array}\right],
$$

where $Q$ sets $a_{i+1}=a_{i}-b_{i}$ and $b_{i+1}=b_{i}$ assuming $a_{i}>b_{i}$, and $S$ is a swap matrix that swaps $a_{i}$ and $b_{i}$ if $a_{i}<b_{i}$. Using these matrices, we implement the algorithm by assigning

$$
\left[\begin{array}{l}
a_{i+1} \\
b_{i+1}
\end{array}\right]=Q\left[\begin{array}{l}
a_{i} \\
b_{i}
\end{array}\right] \text { for } a_{i}>b_{i}, \quad\left[\begin{array}{l}
a_{i+1} \\
b_{i+1}
\end{array}\right]=Q S\left[\begin{array}{l}
a_{i} \\
b_{i}
\end{array}\right] \text { for } a_{i}<b_{i}
$$

The result of this method is a cascade of $Q$ and $S$ matrices. For $(a, b)=(2,3)$, the result is

$$
\left[\begin{array}{l}
1 \\
1
\end{array}\right]=\underbrace{\left[\begin{array}{cc}
1 & -1 \\
0 & 1
\end{array}\right]}_{Q} \underbrace{\left[\begin{array}{ll}
0 & 1 \\
1 & 0
\end{array}\right]}_{S} \underbrace{\left[\begin{array}{cc}
1 & -1 \\
0 & 1
\end{array}\right]}_{Q} \underbrace{\left[\begin{array}{ll}
0 & 1 \\
1 & 0
\end{array}\right]}_{S}\left[\begin{array}{l}
2 \\
3
\end{array}\right]=\underbrace{\left[\begin{array}{cc}
2 & -1 \\
-1 & 1
\end{array}\right]}_{m}\left[\begin{array}{l}
2 \\
3
\end{array}\right] .
$$

Thus we find two solutions $(m, n)$ to the integer equation $2 m+3 n=\operatorname{gcd}(2,3)=1$.

- 6.1: By inspection, find at least one integer pair $(m, n)$ that satisfies $12 m+15 n=3$.

Ans:

- 6.2: Using matrix methods for the Euclidean algorithm, find integer pairs $(m, n)$ that satisfy $12 m+15 n=3$ and $12 m+15 n=0$. Show your work!!!
Ans:
-6.3: Does the equation $12 m+15 n=1$ have integer solutions for $n$ and $m$ ? Why or why not?


## Ans:

## Problem \# 7: Matrix approach:

It can be difficult to keep track of the $a$ 's and $b$ 's when the algorithm has many steps. We need an alternative way to run the Euclidean algorithm using matrix algebra. Matrix methods provide a more transparent approach to the operations on $(a, b)$. Thus the Euclidean algorithm can be classified in terms of standard matrix operations. Write out the indirect matrix approach discussed at the end of Sec. 2.4.3 (Eq. 2.4.3).

Ans:

## Continued fractions

Problem \# 8: Here we explore the continued fraction algorithm (CFA), discussed in Sec. 2.4.4. In its simplest form, the CFA starts with a real number, which we denote as $\alpha \in \mathbb{R}$. Let us work with an irrational real number, $\pi \in \mathbb{I}$, as an example because its CFA representation will be infinitely long. We can represent the CFA coefficients $\alpha$ as a vector of integers $n_{k}, k=1,2, \ldots, \infty$ :

$$
\begin{aligned}
\alpha & =\left[n_{1} ; n_{2}, n_{3}, n_{4}, \ldots\right] \\
& =n_{1}+\frac{1}{n_{2}+\frac{1}{n_{3}+\frac{1}{n_{4}+\cdots}}} .
\end{aligned}
$$

As discussed in Sec. 2.4.3 (p. 27), the CFA is recursive, with three steps per iteration. For $\alpha_{1}=\pi, n_{1}=3, r_{1}=\pi-3$, and $\alpha_{2} \equiv 1 / r_{1}$.

$$
\begin{aligned}
\alpha_{2} & =1 / 0.1416=7.0625 \ldots \\
\alpha_{1}=n_{1}+\frac{1}{\alpha_{2}} & =n_{1}+\frac{1}{n_{2}+\frac{1}{\alpha_{3}}}=\cdots
\end{aligned}
$$

In terms of a Matlab/Octave script,

```
alpha0 = pi;
K=10;
n=zeros(1,K); alpha=zeros(1,K);
alpha(1)=alpha0;
for k=2:K %k=1 to K
n(k)=round(alpha(k-1));
%n(k)=fix(alpha(k-1));
alpha(k)= 1/(alpha(k-1)-n(k));
%disp([fix(k), round(n(k)), alpha(k)]); pause(1)
end
disp([n; alpha]);
%Now compare this to matlab's rat() function
rat(alpha0,1e-20)
```

- 8.1: By hand (you may use Matlab/Octave as a calculator), find the first three values of $n_{k}$ for $\alpha=e^{\pi}$.


## Ans:

- 8.2: For the proceeding question, what is the error (remainder) when you truncate the continued fraction after $n_{1}, \ldots, n_{3}$ ? Give the absolute value of the error and the percentage error relative to the original $\alpha$.
Ans:
- 8.3: Use the Matlab/Octave program provided to find the first 10 values of $n_{k}$ for $\alpha=e^{\pi}$, and verify your result using the Matlab/Octave command rat ().
Ans:
- 8.4: Discuss the similarities and differences between the Euclidean algorithm and the CFA.


## Ans:

Problem \# 9:CFA of ratios of large primes

- 9.1: (4pts) Expand 23/7 as a continued fraction. Express your answer in bracket notation (e.g., $\pi=[3 ., 7,16, \cdots]$ ). Show your work. Ans:
- 9.2:Starting from the primes below $10^{6}$, form the CFA of $\pi_{j} / \pi_{k}$ with $j=78498$ and $k<j$. Ans:
- 9.3: Look at other ratios of prime numbers and look for a pattern in the CFA of the ratios of large primes. What is the most obvious conclusion? Ans:
-9.4: (1pts) Try the Matlab/Octave functions rats (23/7), rats (3.2857), and rats (3.2856). What an you conclude?
Ans:
- 9.5: (2pts) Can $\sqrt{2}$ be represented as a finite continued fraction? Why or why not?

Ans:

- 9.6: (2pts) What is the CFA for $\sqrt{2}-1$ ?

$$
\text { Hint: } \quad \sqrt{2}+1=\frac{1}{\sqrt{2}-1}=[2 ; 2,2,2, \cdots] \text {. }
$$

Ans:

- 9.7: Show that

$$
\frac{1}{1-\sqrt{a}}=a^{\frac{11}{2}}+a^{\frac{9}{2}}+a^{\frac{7}{2}}+a^{\frac{5}{2}}+a^{\frac{3}{2}}+\sqrt{a}+a^{5}+a^{4}+a^{3}+a^{2}+a+1=1-a^{6}
$$

syms a,b
b= taylor(1/( 1-sqrt(a) ))
simplify ((1-sqrt(a))*b) = $1-a^{\wedge} 6$
Use symbolic analysis to show this, then explain. Ans:

### 1.3 Problems NS-3

Topic of this homework: Pythagorean triplets, Pell's equation, Fibonacci sequence

## Pythagorean triplets

Problem \# 1: Euclid's formula for the Pythagorean triplets $a, b$, c is $a=p^{2}-q^{2}, b=2 p q$, and $c=p^{2}+q^{2}$.

- 1.1: What condition(s) must hold for $p$ and $q$ such that $a, b$, and $c$ are always positive and nonzero?


## Ans:

- 1.2: Solve for $p$ and $q$ in terms of $a, b$, and $c$.


## Ans:

Problem \# 2: The ancient Babylonians (ca. 2000 BCE) cryptically recorded ( $a, c$ ) pairs of numbers on a clay tablet, archeologically denoted Plimpton-322 (see 2.8).

- 2.1: Find $p$ and $q$ for the first five pairs of $a$ and $c$ shown here from Plimpton-322.

| $a$ | $c$ |
| :---: | :---: |
| 119 | 169 |
| 3367 | 4825 |
| 4601 | 6649 |
| 12709 | 18541 |
| 65 | 97 |

Find a formula for $a$ in terms of $p$ and $q$.
Ans:

- 2.2: Based on Euclid's formula, show that $c>(a, b)$.

Ans:

- 2.3: What happens when $c=a$ ?


## Ans:

- 2.4: Is $b+c$ a perfect square? Discuss.

Ans:

## Pell's equation:

Problem \# 3: Pell's equation is one of the most historic (i.e., important) equations of Greek number theory because it was used to show that $\sqrt{2} \in \mathbb{I}$. We seek integer solutions of

$$
x^{2}-N y^{2}=1
$$

As shown in Sec. 2.5.2, the solutions $x_{n}, y_{n}$ for the case of $N=2$ are given by the linear $2 \times 2$ matrix recursion

$$
\left[\begin{array}{l}
x_{n+1} \\
y_{n+1}
\end{array}\right]=1 \jmath\left[\begin{array}{ll}
1 & 2 \\
1 & 1
\end{array}\right]\left[\begin{array}{l}
x_{n} \\
y_{n}
\end{array}\right]
$$

with $\left[x_{0}, y_{0}\right]^{T}=[1,0]^{T}$ and $1 \jmath=\sqrt{-1}=e^{j \pi / 2}$. It follows that the general solution to Pell's equation for $N=2$ is

$$
\left[\begin{array}{l}
x_{n} \\
y_{n}
\end{array}\right]=\left(e^{\jmath \pi / 2}\right)^{n}\left[\begin{array}{ll}
1 & 2 \\
1 & 1
\end{array}\right]^{n}\left[\begin{array}{l}
x_{0} \\
y_{0}
\end{array}\right] .
$$

To calculate solutions to Pell's equation using the matrix equation above, we must calculate

$$
A^{n}=e^{\jmath \pi n / 2}\left[\begin{array}{ll}
1 & 2 \\
1 & 1
\end{array}\right]^{n}=e^{\jmath \pi n / 2}\left[\begin{array}{ll}
1 & 2 \\
1 & 1
\end{array}\right]\left[\begin{array}{ll}
1 & 2 \\
1 & 1
\end{array}\right]\left[\begin{array}{ll}
1 & 2 \\
1 & 1
\end{array}\right] \cdots\left[\begin{array}{ll}
1 & 2 \\
1 & 1
\end{array}\right]
$$

which becomes tedious for $n>2$.

- 3.1: Find the companion matrix and thus the matrix $A$ that has the same eigenvalues as Pell's equation. Hint: Use Matlab's function [E, Lambda] = eig (A) to check your results!
Ans:
- 3.2: Solutions to Pell's equation were used by the Pythagoreans to explore the value of $\sqrt{2}$. Explain why Pell's equation is relevant to $\sqrt{2}$.
Ans:
- 3.3: Find the first three values of $\left(x_{n}, y_{n}\right)^{T}$ by hand and show that they satisfy Pell's equation for $N=2$. Ans: By hand, find the eigenvalues $\lambda_{ \pm}$of the $2 \times 2$ Pell's equation matrix

$$
A=\left[\begin{array}{ll}
1 & 2 \\
1 & 1
\end{array}\right]
$$

Ans:

- 3.4: By hand, show that the matrix of eigenvectors, $E$, is

$$
E=\left[\begin{array}{ll}
\vec{e}_{+} & \vec{e}_{-}
\end{array}\right]=\frac{1}{\sqrt{3}}\left[\begin{array}{cc}
-\sqrt{2} & \sqrt{2} \\
1 & 1
\end{array}\right] .
$$

Ans:

- 3.5: Using the eigenvalues and eigenvectors you found for $A$, verify that

$$
E^{-1} A E=\Lambda \equiv\left[\begin{array}{cc}
\lambda_{+} & 0 \\
0 & \lambda_{-}
\end{array}\right]
$$

## Ans:

- 3.6: Once you have diagonalized $A$, use your results for $E$ and $\Lambda$ to solve for the $n=10$ solution $\left(x_{10}, y_{10}\right)^{T}$ to Pell's equation with $N=2$.
Ans:

Problem \# 4: Here we seek the general formula for $x_{n}$. Like Pell's equation, the Fibonacci equation has a recursive eigenanalysis solution. To find it we must recast $x_{n}$ as a $2 \times 2$ matrix relationship and then proceed, as we did for the Pell case.

- 4.1: Show that the Fibonacci sequence $x_{n}=x_{n-1}+x_{n-2}$ may be generated by

$$
\left[\begin{array}{l}
x_{n}  \tag{NS-3.1}\\
y_{n}
\end{array}\right]=\left[\begin{array}{ll}
1 & 1 \\
1 & 0
\end{array}\right]^{n}\left[\begin{array}{l}
x_{0} \\
y_{0}
\end{array}\right], \quad\left[\begin{array}{l}
x_{0} \\
y_{0}
\end{array}\right]=\left[\begin{array}{l}
1 \\
0
\end{array}\right] .
$$

-4.2: What is the relationship between $y_{n}$ and $x_{n}$ ?

## Ans:

- 4.3: Write a Matlab/Octave program to compute $x_{n}$ using the matrix equation above. Test your code using the first few values of the sequence. Using your program, what is $x_{40}$ ? Note: Consider using the eigenanalysis of $A$, described by Eq. 2.5.18 of the text.


## Ans:

-4.4: Using the eigenanalysis of the matrix $A$ (and a lot of algebra), show that it is possible to obtain the general formula for the Fibonacci sequence

$$
\begin{equation*}
x_{n}=\frac{1}{\sqrt{5}}\left[\left(\frac{1+\sqrt{5}}{2}\right)^{n+1}-\left(\frac{1-\sqrt{5}}{2}\right)^{n+1}\right] . \tag{NS-3.2}
\end{equation*}
$$

- 4.5: What are the eigenvalues $\lambda_{ \pm}$of the matrix $A$ ?


## Ans:

- 4.6: How is the formula for $x_{n}$ related to these eigenvalues? Hint: Find the eigenvectors.


## Ans:

-4.7: What happens to each of the two terms

$$
[(1 \pm \sqrt{5}) / 2]^{n+1} ?
$$

Ans:
-4.8: What happens to the ratio $x_{n+1} / x_{n}$ ?
Ans:

Problem \# 5: Replace the Fibonacci sequence with

$$
x_{n}=\frac{x_{n-1}+x_{n-2}}{2},
$$

such that the value $x_{n}$ is the average of the previous two values in the sequence.

- 5.1: What matrix $A$ is used to calculate this sequence?

Ans:
-5.2: Modify your computer program to calculate the new sequence $x_{n}$. What happens as $n \rightarrow \infty$ ?
Ans:

- 5.3: What are the eigenvalues of the modified A? How do they relate to the behavior of $x_{n}$ as $n \rightarrow \infty$ ? Hint: You can expect the closed-form expression for $x_{n}$ to be similar to Eq. NS-3.4.
Ans:

Problem \# 6: Consider the expression

$$
\sum_{1}^{N} f_{n}^{2}=f_{N} f_{N+1}
$$

- 6.1: Find a formula for $f_{n}$ that satisfies this relationship. Hint: It holds for only the Fibonacci recursion formula.
Ans:


## CFA as a matrix recursion

Problem \# 7: The CFA may be writen as a matrix recursion. For this we adopt a special notation, unlike other matrix notations, ${ }^{2}$ with $k \in \mathbb{N}$ :

$$
\left[\begin{array}{l}
n \\
x
\end{array}\right]_{k+1}=\left[\begin{array}{cc}
0 & \left\lfloor x_{k}\right\rfloor \\
0 & \frac{1}{x_{k}-\left\lfloor x_{k}\right\rfloor}
\end{array}\right]\left[\begin{array}{l}
n \\
x
\end{array}\right]_{k} .
$$

This equation says that $n_{k+1}=\left\lfloor x_{k}\right\rfloor$ and $x_{k+1}=1 /\left(x_{k}-\left\lfloor x_{k}\right\rfloor\right)$. It does not mean that $n_{k+1}=\left\lfloor x_{k}\right\rfloor x_{k}$, as would be implied by standard matrix notation. The lower equation says that $r_{k}=x_{k}-\left\lfloor x_{k}\right\rfloor$ is the remainder-namely, $x_{k}=$ $\lfloor x-k\rfloor+r_{k}$ (Octave/Matlab's rem ( $\mathrm{x}, \mathrm{floor}(\mathrm{x})$ ) function), also known as mod ( $\mathrm{x}, \mathrm{y}$ ).

- 7.1: Start with $n_{0}=0 \in \mathbb{N}, x_{0} \in \mathbb{I}, n_{1}=\left\lfloor x_{0}\right\rfloor \in \mathbb{N}, r_{1}=x-\lfloor x\rfloor \in \mathbb{I}$, and $x_{1}=1 / r_{1} \in \mathbb{I}, r_{n} \neq 0$. For $k=1$ this generates on the left the next CFA parameter $n_{2}=\left\lfloor x_{1}\right\rfloor$ and $x_{2}=1 / r_{2}=1 /\left(x_{0}-\left\lfloor x_{0}\right\rfloor\right)$ from $n_{0}$ and $x_{0}$. Find $[n, x]_{k+1}^{T}$ for $k=2,3,4,5$.
Ans:

[^1]
## Chapter 2

## Algebraic Equations

### 2.1 Problems AE-1

Topics of this homework: Fundamental theorem of algebra, polynomials, analytic functions and their inverse, convolution, Newton's root finding method, Riemann zeta function. Deliverables: Answers to problems

Note: The term analytic is used in two different ways. (1) An analytic function is a function that may be expressed as a locally convergent power series; (2) analytic geometry refers to geometry using a coordinate system.

## Polynomials and the fundamental theorem of algebra (FTA)

Problem \# 1: A polynomial of degree $N$ is defined as

$$
P_{N}(x)=a_{0}+a_{1} x+a_{2} x^{2}+\cdots+a_{N} x^{N} .
$$

- 1.1: How many coefficients $a_{n}$ does a polynomial of degree $N$ have?


## Ans:

- 1.2: How many roots does $P_{N}(x)$ have?


## Ans:

Problem \# 2: The fundamental theorem of algebra (FTA)

## Ans:

- 2.2: Using the FTA, prove your answer to question 1.2. Hint: Apply the FTA to prove how many roots a polynomial $P_{N}(x)$ of order $N$ has.
Ans:

Problem \# 3: Consider the polynomial function $P_{2}(x)=1+x^{2}$ of degree $N=2$ and the related function $F(x)=1 / P_{2}(x)$. What are the roots (e.g., zeros) $x_{ \pm}$of $P_{2}(x)$ ? Hint: Complete the square on the polynomial $P_{2}(x)=1+x^{2}$ of degree 2 , and find the roots.
Ans:

Problem \# 4: $F(x)$ may be expressed as $\left(A, B, x_{ \pm} \in \mathbb{C}\right)$

$$
\begin{equation*}
F(x)=\frac{A}{x-x_{+}}+\frac{B}{x-x_{-}} \tag{AE-1.1}
\end{equation*}
$$

where $x_{ \pm}$are the roots (zeros) of $P_{2}(x)$, which become the poles of $F(x) ; A$ and $B$ are the residues. The expression for $F(x)$ is sometimes called a partial fraction expansion or residue expansion, and it appears in many engineering applications.

- 4.1: Find $A, B \in \mathbb{C}$ in terms of the roots $x_{ \pm}$of $P_{2}(x)$.


## Ans:

-4.2: Verify your answers for $A$ and $B$ by showing that this expression for $F(x)$ is indeed equal to $1 / P_{2}(x)$.
Ans:

- 4.3: Give the values of the poles and zeros of $P_{2}(x)$.


## Ans:

- 4.4: Give the values of the poles and zeros of $F(x)=1 / P_{2}(x)$.

Ans:

### 2.1.1 Analytic functions

Overview: Analytic functions are defined by infinite (power) series. The function $f(x)$ is said to be analytic at any value of constant $x=x_{o}$, where there exists a convergent power series

$$
P(x)=\sum_{n=0}^{\infty} a_{n}\left(x-x_{o}\right)^{n}
$$

such that $P\left(x_{o}\right)=f\left(x_{o}\right)$. The point $x=x_{o}$ is called the expansion point. The region around $x_{o}$ such that $\left|x-x_{o}\right|<1$ is called the radius of convergence, or region of convergence (RoC). The local power series for $f(x)$ about $x=x_{o}$ is defined by the Taylor series:

$$
\begin{aligned}
f(x) & \approx f\left(x_{o}\right)+\left.\frac{d f}{d x}\right|_{x=x_{o}}\left(x-x_{o}\right)+\left.\frac{1}{2!} \frac{d^{2} f}{d x^{2}}\right|_{x=x_{o}}\left(x-x_{o}\right)^{2}+\cdots \\
& =\left.\sum_{n=0}^{\infty} \frac{1}{n!} \frac{d^{n}}{d x^{n}} f(x)\right|_{x=x_{o}}\left(x-x_{o}\right)^{n}
\end{aligned}
$$

Two classic examples are the geometric series ${ }^{1}$ where $a_{n}=1$,

$$
\begin{equation*}
\frac{1}{1-x}=1+x+x^{2}+x^{3}+\cdots=\sum_{n=0}^{\infty} x^{n} \tag{AE-1.2}
\end{equation*}
$$

and the exponential function where $a_{n}=1 / n$ !, Eq. 3.2 .11 (p. 70). The coefficients for both series may be derived from the Taylor formula.

## Problem \# 5: The geometric series

- 5.1: What is the region of convergence $($ RoC $)$ for the power series Eq. AE-1.2 of $1 /(1-x)$ given above-for example, where does the power series $P(x)$ converge to the function value $f(x)$ ? State your answer as a condition on $x$. Hint: What happens to the power series when $x>1$ ?
Ans:
- 5.2: In terms of the pole, what is the RoC for the geometric series in Eq. AE-1.2?


## Ans:

[^2]- 5.3: How does the RoC relate to the location of the pole of $1 /(1-x)$ ?


## Ans:

- 5.4: Where are the zeros, if any, in Eq. AE-1.2?


## Ans:

- 5.5: Assuming $x$ is in the RoC, prove that the geometric series correctly represents $1 /(1-x)$ by multiplying both sides of Eq. AE-1.2 by $(1-x)$.


## Ans:

Problem \# 6: Use the geometric series to study the degree $N$ polynomial. It is very important to note that all the coefficients $c_{n}$ of this polynomial are 1.

$$
\begin{equation*}
P_{N}(x)=1+x+x^{2}+\cdots+x^{N}=\sum_{n=0}^{N} x^{n} . \tag{AE-1.3}
\end{equation*}
$$

- 6.1: Prove that

$$
\begin{equation*}
P_{N}(x)=\frac{1-x^{N+1}}{1-x} . \tag{AE-1.4}
\end{equation*}
$$

Ans:

## Ans:

- 6.3: What is the RoC for Eq. AE-1.4?


## Ans:

- 6.4: How many poles does $P_{N}(x)(E q . A E-1.3)$ have? Where are they?


## Ans:

-6.5: How many zeros does $P_{N}(x)(E q . A E-1.4)$ have? State where are they in the complex plane. Ans:

- 6.6: Explain why Eqs. AE-1.3 and AE-1.4 have different numbers of poles and zeros.


## Ans:

- 6.7: Is the function $1 /(1-x)$ analytic outside of the RoC?


## Ans:

- 6.8: Extra credit. Evaluate $P_{N}(x)$ at $x=0$ and $x=0.9$ for the case of $N=100$, and compare the result to that from Matlab.

```
%sum the geometric series and P_100(0.9)
clear all;close all;format long
N=100; x=0.9; S=0;
for n=0:N
S=S+x^n
```

```
end
P100=(1-x^ (N+1))/(1-x);
disp(sprintf('S= %g, P100= %g, error= %g',S,P100, S-P100))
```

Ans:

Problem \# 7: The exponential series

- 7.1: What is the RoC for the exponential series Eq. 3.2.11?


## Ans:

- 7.2: Let $x=\jmath$ in Eq. 3.2.11, and write out the series expansion of $e^{x}$ in terms of its real and imaginary parts.


## Ans:

- 7.3: Let $x=\jmath \theta$ in Eq. 3.2.11, and write out the series expansion of $e^{x}$ in terms of its real and imaginary parts. How does your result relate to Euler's identity $\left(e^{\jmath \theta}=\cos (\theta)+\jmath \sin (\theta)\right)$ ?


## Ans:

### 2.1.2 Inverse analytic functions and composition

Overview: It may be surprising, but every analytic function has an inverse function. Starting from the function $(x, y \in \mathbb{C})$

$$
y(x)=\frac{1}{1-x}
$$

the inverse is

$$
x=\frac{y-1}{y}=1-\frac{1}{y} .
$$

Problem \# 8: Consider the inverse function described above
$-8.1:$ Where are the poles and zeros of $x(y)$ ?

## Ans:

-8.2: Where (for what condition on $y$ ) is $x(y)$ analytic?
Ans:

Problem \# 9 Consider the exponential function $z(x)=e^{x}(x, z \in \mathbb{C})$.
-9.1: Find the inverse $x(z)$.

## Ans:

-9.2: Where are the poles and zeros of $x(z)$ ?

## Ans:

Problem \# 10: Composition.

- 10.1: If $y(s)=1 /(1-s)$ and $z(s)=e^{s}$, compose these two functions to obtain $(y \circ z)(s)$. Give the expression for $(y \circ z)(s)=y(z(s))$. Ans:
$-10.2:$ Where are the poles and zeros of $(y \circ z)(s)$ ?


## Ans:

- 10.3: Where (for what condition on $x$ ) is $(y \circ z)(x)$ analytic?


## Ans:

## Eigen-analysis

Problem \# 11: (4 pts) The vectorized eigen-equation for a matrix $\mathfrak{A}$ is

$$
\begin{equation*}
\boldsymbol{A} \boldsymbol{E}=\boldsymbol{E} \boldsymbol{\Lambda} \tag{AE-1.5}
\end{equation*}
$$

- 11.1: (4pt) Provide a formula for $\boldsymbol{A}^{3}$ in terms of the eigenvector $\boldsymbol{E}$ and eigenvalue $\boldsymbol{\Lambda}$ matricies. Ans:
- 11.2: (4 pts) Find the eigenvalues of the matrix, and find the roots, by completing the square, where $a, b, c, d \in \mathbb{C}$, and

$$
\boldsymbol{A}=\left[\begin{array}{ll}
a & b \\
c & d
\end{array}\right] .
$$

Ans:

## (4 pts) Convolution

Multiplying two short or simple polynomials is not demanding. However, if the polynomials have many terms, it can become tedious. For example, multiplying two 10th-degree polynomials is not something one would like to do every day.

An alternative is a method called convolution. The inverse of convolution is called deconvolution, which is equivalent to long-division of polynomials, also known as factoring polynomials (Sec. 3.4.1, p. 109-111). Newton's method is a reliable and accurately algorithm to extract roots from polynomials using term by term deconvolution. When the roots are well approximated by fractional numbers, the method is accurate to within computational accuracy. For example, if the root is $\pi \approx \hat{\pi}_{19} \equiv 817696623 / 260280919 \in \mathbb{F}$, as given by rats $(\mathrm{pi}, 19) . \hat{\pi}_{19}$ is the 64 bit machine's internal representation of $\pi$ since $\pi-\hat{\pi}_{19}=0$ (See text Fig. 2.6, p. 48).

Problem \# 12: (4 pts) Convolution of sequences. Practice convolution (by hand!!) using a few simple examples. Manually evaluate the following convolutions. Show your work!

- 12.1: (2 pts) Multiplying two polynomials is the same as convolving their coefficients.

Given

$$
\begin{aligned}
f(x) & =x^{3}+3 x^{2}+3 x+1 \leftrightarrow[1 ; 3,3,1] \\
g(x) & =x^{3}+2 x^{2}+x+2 \leftrightarrow[1 ; 2,1,2] .
\end{aligned}
$$

show that

$$
f(x) g(x)=x^{6}+5 x^{5}+10 x^{4}+12 x^{3}+11 x^{2}+7 x+2 \leftrightarrow[1 ; 3,3,1] \star[1 ; 2,1,2] .
$$

Ans:

$$
-12.2:(1 \mathrm{pts})[1 ;-1] \star[1 ; 2,4,7,0]
$$

## Ans:

$$
-12.3:(1 \mathrm{pts})[1 ; 2,1] \star[1 ;-1]
$$

## Ans:

## Newton's root-finding method

Problem \# 13: Use Newton's iteration to find the roots of the polynomial

$$
P_{3}(x)=1-x^{3} .
$$

- 13.1: Draw a graph describing the first step of the iteration starting with $x_{0}=(1 / 2,0)$.


## Ans:

- 13.2: Calculate $x_{1}$ and $x_{2}$. What number is the algorithm approaching?

Ans:

- 13.3: Does Newton's method work for $P_{2}(x)=1+x^{2}$ ? If so, why? Hint: What are the roots in this case?
Ans:

Problem \# 14: In this problem we consider the case of fractional roots, and take advantage of this fact during the itteration. Given that the roots are integers, composed of primes, we may uniquely identify the primes by factoring the numerator and denominator of the rational approximation of the root. The method is:

1. Start the Newton itteration

$$
s_{n+1}=s_{n}-\frac{M\left(s_{n}\right)}{M^{\prime}\left(s_{n}\right)}
$$

2. Apply the CFA to the next output rats $\left(s_{n+1}\right)$
3. Factor the Num and Dem of the CFA
4. Terminate when the factors converge

Using this method, show that we can find either the best possible fractional approximation to the roots (or even the exact roots, when the answer is within machine accuracy).

- 14.1: Find the roots of a Monic having coefficents $m_{k} \in \mathbb{F}$.

Let

$$
M_{3}(x)=(x-254 / 17)(x-2047 / 13)(x-17 / 13)
$$

In this case the root vector $R$ becomes

$$
R=[14.9412,157.4615,1.3077] .
$$

Verify that rats (M) returns the rational set of roots. Ans:

### 2.2 Problems AE-2

## Topics of this homework:

Linear vs nonlinear systems of equations, Euclid's formula, Gaussian elimination, matrix permutations, Ohm's law, twoport networks,

Deliverables: Answers to problems

## Gaussian elimination

Problem \# 1: Gaussian elimination

- 1.1: Find the inverse of

$$
A=\left[\begin{array}{ll}
1 & 2 \\
4 & 3
\end{array}\right]
$$

Ans:

- 1.2: Verify that $A^{-1} A=A A^{-1}=\left[\begin{array}{ll}1 & 0 \\ 0 & 1\end{array}\right]$.

Ans:

Problem \# 2: Find the solution to the following $3 \times 3$ matrix equation $A x=b$ by GE. Show your intermediate steps. You can check your work at each step using Octave/Matlab.

$$
\left[\begin{array}{ccc}
1 & 1 & -1 \\
3 & 1 & 1 \\
1 & -1 & 4
\end{array}\right]\left[\begin{array}{l}
x_{1} \\
x_{2} \\
x_{3}
\end{array}\right]=\left[\begin{array}{l}
1 \\
9 \\
8
\end{array}\right] .
$$

- 2.1 Show (i.e., verify) that the first GE matrix $G_{1}$, which zeros out all entries in the first column is given by

$$
G_{1}=\left[\begin{array}{ccc}
1 & 0 & 0 \\
-3 & 1 & 0 \\
-1 & 0 & 1
\end{array}\right]
$$

Identify the elementary row operations that this matrix performs. Ans:

- 2.2 Find a second GE matrix, $G_{2}$, to put $G_{1} A$ in upper triangular form. Identify the elementary row operations that this matrix performs.


## Ans:

- 2.3 Find a third GE matrix $G_{3}$ that scales each row so that its leading term is 1. Identify the elementary row operations that this matrix performs.


## Ans:

-2.4: Find the last GE matrix, $G_{4}$, which subtracts a scaled version of row 3 from row 2, and scaled versions of rows 2 and 3 from row 1, such that you are left with the identity matrix ( $G_{4} G_{3} G_{2} G_{1} A=I$ ).

## Ans:

-2.5: Solve for $\left\{x_{1}, x_{2}, x_{3}\right\}^{T}$ using the augmented matrix format $G_{4} G_{3} G_{2} G_{1}\{A \mid b\}$ (where $\{A \mid b\}$ is the augmented matrix). Note that if you've performed the preceding steps correctly, $x=G_{4} G_{3} G_{2} G_{1} b$. Ans:

- 2.6: Find the pivot matrix $G$ that rescales the second row of the augmented matrix $A \mid b$ by $1 / 3$. Ans:


## Two linear equations

Problem \# 3 In this problem we transition from a general pair of equations

$$
\begin{aligned}
& f(x, y)=0 \\
& g(x, y)=0
\end{aligned}
$$

to the important case of two linear equations

$$
\begin{aligned}
& y=a x+b \\
& y=\alpha x+\beta .
\end{aligned}
$$

Note that to help keep track of the variables, roman coefficients $(a, b)$ are used for the first equation and Greek $(\alpha, \beta)$ for the second.

- 3.1: What does it mean, graphically, if these two linear equations have (1) a unique solution, (2) a nonunique solution, or (3) no solution?
Ans:
- 3.2: Assuming the two equations have a unique solution, find the solution for $x$ and $y$.


## Ans:

- 3.3: When will this solution fail to exist (for what conditions on $a, b, \alpha$, and $\beta$ )?


## Ans:

-3.4: Write the equations as a $2 \times 2$ matrix equation of the form $A \vec{x}=\vec{b}$, where $\vec{x}=\{x, y\}^{T}$.

## Ans:

- 3.5: Find the inverse of the $2 \times 2$ matrix, and solve the matrix equation for $x$ and $y$.


## Ans:

- 3.6: Discuss the properties of the determinant of the matrix ( $\Delta$ ) in terms of the slopes of the two equations ( $a$ and $\alpha$ ).


## Ans:

Problem \# 4: The application of linear functional relationships between two variables
We use $2 \times 2$ matrices to describe two-port networks, as discussed in Sec. 3.8 (p. 108). Transmission lines are a great example: Both voltage and current must be tracked as they travel along the line. Figure 3.8 (p. 112) shows an example segment of a transmission line.

Suppose you are given the following pair of linear relationships between the input (source) variables $V_{1}$ and $I_{1}$ and the output (load) variables $V_{2}$ and $I_{2}$ of the transmission line:

$$
\left[\begin{array}{c}
V_{1} \\
I_{1}
\end{array}\right]=\left[\begin{array}{cc}
3 & 1 \\
1 & -1
\end{array}\right]\left[\begin{array}{l}
V_{2} \\
I_{2}
\end{array}\right] .
$$

-4.1: Let the output (the load) be $V_{2}=1$ and $I_{2}=2$ (i.e., $V_{2} / I_{2}=1 / 2\{\Omega\}$ ). Find the input voltage and current, $V_{1}$ and $I_{1}$.

## Ans:

-4.2: Let the input (source) be $V_{1}=1$ and $I_{1}=2$. Find the output voltage and current, $V_{2}$ and $I_{2}$.

## Ans:

## Integer equations: applications and solutions

Any equation for which we seek only integer solutions is called a Diophantine equation.

## Problem \# 5: A practical example of using a Diophantine equation:

"A merchant had a 40-pound weight that broke into 4 pieces. When the pieces were weighed, it was found that each piece was a whole number of pounds and that the four pieces could be used to weigh every integral weight between 1 and 40 pounds. What were the weights of the pieces?" - Bachet de Bèziriac (1623) ${ }^{2}$.

Here, weighing is performed using a balance scale that has two pans, with weights on either pan. Thus, given weights of 1 and 3 pounds, one can weigh a 2-pound weight by putting the 1-pound weight in the same pan with the 2-pound weight, and the 3 -pound weight in the other pan. Then the scale will be balanced. A solution to the four weights for Bachet's problem is $1+3+9+27=40$ pounds.

[^3]

- 5.1: Show how the combination of $1-, 3-, 9-$, and 27 -pound weights can be used to weigh $1,2,3, \ldots, 8,28$, and 40 pounds of milk (or something else, such as flour). Assuming that the milk is in the left pan, provide the position of the weights using a negative sign - to indicate the left pan and a positive sign + to indicate the right pan. For example, if the left pan has 1 pound of milk, then 1 pound of milk in the right pan, +1 , will balance the scales.

Hint: It is helpful to write the answer in matrix form. Set the vector of values to be weighed equal to a matrix indicating the pan assignments, multiplied by a vector of the weights $[1,3,9,27]^{T}$. The pan assignments matrix should contain only the values -1 (left pan), +1 (right pan), and 0 (leave out). You can indicate these using,-+ , and blanks.

Ans:

## Vector algebra in $\mathbb{R}^{3}$

Definitions of the scalar (also called a dot product) $\boldsymbol{A} \cdot \boldsymbol{B}$, cross $\boldsymbol{A} \times \boldsymbol{B}$ and triple product $\boldsymbol{A} \cdot(\boldsymbol{B} \times \boldsymbol{C})$, may be found in Appendix A (p.215), where $\boldsymbol{A}, \boldsymbol{B}, \boldsymbol{C}$ in $\mathbb{R}^{3} \subset \mathbb{C}^{3}$ ), as shown in Fig. 3.4, p. 87. A fourth "double-cross" ( is. ${ }^{3}$

$$
\boldsymbol{A} \times(\boldsymbol{B} \times \boldsymbol{C})=\alpha_{o} \boldsymbol{B}-\beta_{o} \boldsymbol{C}
$$

where $\alpha_{o}=\boldsymbol{A} \cdot \boldsymbol{C}$ and $\beta_{o}=\boldsymbol{A} \cdot \boldsymbol{B}($ Note: $\boldsymbol{A} \times(\boldsymbol{B} \times \boldsymbol{C}) \neq(\boldsymbol{A} \times \boldsymbol{B}) \times \boldsymbol{C})$.
Problem \# 6: Scalar product A • B

- 6.1: If $\boldsymbol{A}=a_{x} \hat{\mathbf{x}}+a_{y} \hat{\mathbf{y}}+a_{z} \hat{\mathbf{z}}$ and $\boldsymbol{B}=b_{x} \hat{\mathbf{x}}+b_{y} \hat{\mathbf{y}}+b_{z} \hat{\mathbf{z}}$, write out the definition of $\mathbf{A} \cdot \mathbf{B}$.


## Ans:

- 6.2: The dot product is often defined as $\|\boldsymbol{A}\|\|\boldsymbol{B}\| \cos (\theta)$, where $\|\boldsymbol{A}\|=\sqrt{\boldsymbol{A} \cdot \boldsymbol{A}}$ and $\theta$ is the angle between $\boldsymbol{A}, \boldsymbol{B}$. If $\|\boldsymbol{A}\|=1$, describe how the dot product relates to the vector $\mathbf{B}$.
Ans:

[^4]Problem \# 7: Vector (cross) product $\mathrm{A} \times \mathrm{B}$

- 7.1: If $\boldsymbol{A}=a_{x} \hat{\mathbf{x}}+a_{y} \hat{\mathbf{y}}+a_{z} \hat{\mathbf{z}}$ and $\boldsymbol{B}=b_{x} \hat{\mathbf{x}}+b_{y} \hat{\mathbf{y}}+b_{z} \hat{\mathbf{z}}$, write out the definition of $\mathbf{A} \times \mathbf{B}$.

Ans:

- 7.2: Show that the cross product is equal to the area of the parallelogram formed by $\boldsymbol{A}, \boldsymbol{B}$, namely $\|\boldsymbol{A}\|\|\boldsymbol{B}\| \sin (\theta)$, where $\|\boldsymbol{A}\|=\sqrt{\boldsymbol{A} \cdot \boldsymbol{A}}$ and $\theta$ is the angle between $\boldsymbol{A}$ and $\boldsymbol{B}$.
Ans:

Problem \# 8: Triple product $\boldsymbol{A} \cdot(\boldsymbol{B} \times \boldsymbol{C})$
Let $\boldsymbol{A}=\left[a_{1}, a_{2}, a_{3}\right]^{T}, \boldsymbol{B}=\left[b_{1}, b_{2}, b_{3}\right]^{T}, \boldsymbol{C}=\left[c_{1}, c_{2}, c_{3}\right]^{T}$ be three vectors in $\mathbb{R}^{3}$.

- 8.1: Describe why $|\boldsymbol{A} \cdot(\boldsymbol{B} \times \boldsymbol{C})|$ is the volume of parallelepiped generated by $\boldsymbol{A}, \boldsymbol{B}$, and $\boldsymbol{C}$.

Ans:

- 8.2: Explain why three vectors $\mathbf{A}, \mathrm{B}, \mathrm{C}$ are in one plane if and only if the triple product $\mathrm{A} \cdot(\mathrm{B} \times \mathrm{C})=0$.
Ans:

Problem \# 9: Given two vectors $\boldsymbol{A}, \boldsymbol{B}$ in the $\hat{\mathbf{x}}, \hat{\mathbf{y}}$ plane shown in Fig.3.4 with $\boldsymbol{B}=\hat{\mathbf{y}}(i . e .,\|\boldsymbol{B}\|=1$ ).

- 9.1: Show that $\boldsymbol{A}$ may be split into two orthogonal parts, one in the direction of $\boldsymbol{B}$ and the other perpendicular ( $\perp$ ) to $\boldsymbol{B}$. Hint: Express the vector products of $\boldsymbol{A}$ and $\boldsymbol{B}$ (dot and cross) in polar coordinates (Greenberg, 1988).

$$
\begin{aligned}
\boldsymbol{A} & =(\boldsymbol{A} \cdot \boldsymbol{B}) \boldsymbol{B}+\boldsymbol{B} \times(\boldsymbol{A} \times \boldsymbol{B}) \\
& =\boldsymbol{A}_{\|}+\boldsymbol{A}_{\perp}
\end{aligned}
$$

## Ans:

## Ohm's Law

In general, impedance is defined as the ratio of a force to a flow. For electrical circuits, the voltage is the force and the current is the flow. Ohm's law states that the voltage across and the current through a circuit element are related by the impedance of that element (which may be a function of frequency). For resistors, the voltage over the current is called the resistance and is a constant (e.g., the simplest case is $V / I=R$ ). For inductors and capacitors, the voltage over the current is a frequency-dependent impedance (e.g., $V / I=Z(s)$, where $s$ is the complex frequency $s \in \mathbb{C}$ ).

As shown in Table 3.2 (p. 111), the impedance concept also holds in mechanics and acoustics. In mechanics, the force is equal to the mechanical force on an element (e.g., a mass, dashpot, or spring) and the flow is the velocity. In acoustics, the force is pressure and the flow is the volume velocity or particle velocity of air molecules.

| Case | Force | Flow | Impedance | units |
| :--- | :---: | :---: | :---: | :---: |
| Electrical | voltage (V) | current (I) | $Z$ | Ohms $[\Omega]$ |
| Mechanics | force (F) | velocity (V) | $Z$ | Mechanical Ohms $[\Omega]$ |
| Acoustics | pressure (P) | particle velocity (U) | $Z$ | Acoustic Ohms $[\Omega]$ |
| Thermal | temperature (T) | heat-flux (J) | $Z$ | Thermal Ohms $[\Omega]$ |

Problem \# 10: The resistance of an incandescent (filament) lightbulb, measured cold, is about 100 ohms. As the bulb lights up, the resistance of the metal filament increases.
Ohm's law says that the current

$$
\frac{V}{I}=R(T),
$$

where $T$ is the temperature. In the United States, the voltage is 120 volts (RMS) at $60[\mathrm{~Hz}]$. Find the current when the light is first switched on. Ans:

Problem \# 11: The power in watts is the product of the force and the flow. What is the power of the lightbulb of Problem 10?

## Ans:

Problem \# 12: State the impedance Z(s) of each of the following circuit elements: (1) a resistor with resistance $R$, (2) an inductor with inductance $L$, and (3) a capacitor with capacitance $C$.

Ans:

Problem \# 13: Describe the temperature as a function of time when 1 kg of water, 1 kg of ice and 1 kgm of steam are place in an insulated chamber.

## Ans:

Problem \# 14: Consider what happens at the triple point of water. As water freezes or thaws, the temperature remains constant at $0\left(C^{\circ}\right)$. Once all the water is frozen and more heat is removed, the temperature drops below $0^{\circ}$. As heat is added, water thaws but the temperature remains at $0^{\circ}$. Once all the ice has melted, what is the temperature as more heat is added?
Model the triple point using a Zener diode, a resistor, and a capacitor. A Zener diode holds the voltage constant independent of current. For the case of water's triple point, the voltage represents the temperature of water at the triple point, clamped at $0\left[\mathrm{C}^{\circ}\right]$. The current represents the heat flux. The latent heat of water at the triple point is $32 \mathrm{Cal} / \mathrm{gm}$. Thus as the temperature rises from below freezing, the water is clamped at $0^{\circ}$ once the triple point is reached. At that point, adding more heat flux has no effect on the temperature until all the ice melts. Once the ice has melted, the temperature again begins to rise until it hits the boiling point, where it again stays at $100^{\circ}$ until all the water has evaporated. Ans:

## Nonlinear (quadratic) to linear equations

In the following problems we deal with algebraic equations in more than one variable that are not linear equations. For example, the circle $x^{2}+y^{2}=1$ may be solved for $y(x)= \pm \sqrt{1-x^{2}}$. If we let $z_{+}=x+y \jmath=x+\jmath \sqrt{1-x^{2}}=e^{\theta \jmath}$, we obtain the equation for half a circle $(y>0)$. The entire circle is described by the magnitude of $z$ as $|z|^{2}=(x+y \jmath)(x-y \jmath)=$ 1.

Problem \# 15: Give the curve defined by the equation:

$$
x^{2}+x y+y^{2}=1
$$

- 15.1: Find the function $y(x)$.


## Ans:

- 15.2: Using Matlab/Octave, plot $y(x)$ and describe the graph.

Ans:

- 15.3: What is the name of this curve?


## Ans:

- 15.4: Find the solution (in $x, p$, and $q$ ) to these equations:

$$
\begin{gathered}
x+y=p \\
x y=q .
\end{gathered}
$$

Ans:

- 15.5: Find an equation that is linear in $y$ starting from equations that are quadratic (second-degree) in the two unknowns $x$ and $y$ :

$$
\begin{align*}
x^{2}+x y+y^{2} & =1  \tag{AE-2.1}\\
4 x^{2}+3 x y+2 y^{2} & =3 \tag{AE-2.2}
\end{align*}
$$

Ans:

- 15.6: Compose the following two quadratic equations and describe the results.

$$
\begin{aligned}
x^{2}+x y+y^{2} & =1 \\
2 x^{2}+x y & =1
\end{aligned}
$$

## Ans:

## Nonlinear intersection in analytic geometry

Euclid's formula for Pythagorean triplets (Eq. 2.5.6, p. 40) can be derived by intersecting a circle and a secant line. Consider the nonlinear equation of a unit circle having radius 1 , centered at $(x, y)=(0,0)$,

$$
x^{2}+y^{2}=1,
$$

and the secant line through $(-1,0)$,

$$
y=t(x+1)
$$

a linear equation having slope $t$ and intercept $x=-1$. If the slope $0<t<1$, the line intersects the circle at a second point $(a, b)$ in the positive $x, y$ quadrant. The goal is to find $a, b \in \mathbb{N}$ and then show that $c^{2}=a^{2}+b^{2}$. Since the construction gives a right triangle with short sides $a, b \in \mathbb{N}$, then it follows that $c \in \mathbb{N}$.


Figure 2.1: Derivation of Euclid's formula for the Pythagorean triplets $(P T)[a, b, c]$, based on a composition of a line, having a rational slope $t=p / q \in \mathbb{F}$, and a circle $c^{2}=a^{2}+b^{2},[a, b, c] \in \mathbb{N}$. This analysis is attributed to Diophantus (Di•o•phan' tus) ( 250 CE), and today such equations are called Diophantine (Di.o•phan'tine) equations. PTs have applications in architecture and scheduling, and many other practical problems. Most interesting is their relation to Rydberg's formula for the eigenstates of the hydrogen atom (Appendix H).

Problem \# 16: Derive Euclid's formula

- 16.1: Draw the circle and the line, given a positive slope $0<t<1$.


## Ans:

Problem \# 17: Substitute $y=t(x+1)$ (the line equation) into the equation for the circle, and solve for $x(t)$.
Hint: Because the line intersects the circle at two points, you will get two solutions for $x$. One of these solutions is the trivial solution $x=-1$. Ans:

- 17.1: Substitute the $x(t)$ you found back into the line equation, and solve for $y(t)$.


## Ans:

- 17.2: Let $t=q / p$ be a rational number, where $p$ and $q$ are integers. Find $x(p, q)$ and $y(p, q)$.


## Ans:

- 17.3: Substitute $x(p, q)$ and $y(p, q)$ into the equation for the circle, and show how Euclid's formula for the Pythagorean triples is generated.


## Ans:

For full points you must show that you understand the argument. Explain the meaning of the comment "magic happens" when $t^{4}$ cancels.

### 2.3 Problems AE-3

## Topics of this homework:

Visualizing complex functions, bilinear/Möbius transformation, Riemann sphere.
Deliverables: Answers to problems

## Two-port network analysis

Problem \# 1: Perform an analysis of electrical two-port networks, shown in Fig. 3.6 (page 144). This can be a mechanical system if the capacitors are taken to be springs and inductors taken as mass, as in the suspension of the wheels of a car. In an acoustical circuit, the low-pass filter could be a car muffler. While the physical representations will be different, the equations and the analysis are exactly the same.
The definition of the ABCD transmission matrix $(\mathcal{T})$ is

$$
\left[\begin{array}{c}
V_{1}  \tag{AE-3.1}\\
I_{1}
\end{array}\right]=\left[\begin{array}{ll}
\mathcal{A} & \mathcal{B} \\
\mathcal{C} & \mathcal{D}
\end{array}\right]\left[\begin{array}{c}
V_{2} \\
-I_{2}
\end{array}\right] .
$$

The impedance matrix, where the determinant $\Delta_{\mathcal{T}}=A D-B C$, is given by

$$
\left[\begin{array}{l}
V_{1}  \tag{AE-3.2}\\
V_{2}
\end{array}\right]=\frac{1}{\mathcal{C}}\left[\begin{array}{cc}
\mathcal{A} & \Delta_{\mathcal{T}} \\
1 & \mathcal{D}
\end{array}\right]\left[\begin{array}{l}
I_{1} \\
I_{2}
\end{array}\right] .
$$

- 1.1: Derive the formula for the impedance matrix (Eq. AE-3.2) given the transmission matrix definition (Eq. AE-3.1). Show your work.


## Ans:

Problem \# 2: Consider a single circuit element with impedance $Z(s)$.

- 2.1: What is the $A B C D$ matrix for this element if it is in series?

Ans:

- 2.2: What is the ABCD matrix for this element if it is in shunt?


## Ans:

Problem \# 3: Find the ABCD matrix for each of the circuits of Fig. 3.6.
For each circuit, (i) show the cascade of transmission matrices in terms of the complex frequency $s \in \mathbb{C}$, then (ii) substitute $s=1 \jmath$ and calculate the total transmission matrix at this single frequency.

- 3.1: Left circuit (let $R_{1}=R_{2}=10$ kilo-ohms and $C=10$ nano-farads)


## Ans:

- 3.2: Right circuit (use L and C values given in the figure), where the pressure $P$ is analogous to the voltage $V$, and the velocity $U$ is analogous to the current $I$.


## Ans:

- 3.3: Convert both transmission (ABCD) matrices to impedance matrices using Eq. AE-3.2. Do this for the specific frequency $s=1$ 〕 as in the previous part (feel free to use Matlab/Octave for your computation).


## Ans:

- 3.4: Right circuit: Repeat the analysis as in question 3.3.


## Ans:

## Algebra

Problem \# 4: Fundamental theorem of algebra (FTA).

- 4.1: State the fundamental theorem of algebra (FTA).


## Ans:

## ( 13 pts ) Algebra with complex variables

Problem \# 5: (7 pts) Order and complex numbers:
One can always say that $3<4$-namely, that real numbers have order. One way to view this is to take the difference and compare it to zero, as in $4-3>0$. Here we will explore how complex variables may be ordered. In the following define $\{x, y\} \in \mathbb{R}$ and complex variable $z=x+y \jmath \in \mathbb{C}$.

- 5.1: Explain the meaning of $\left|z_{1}\right|>\left|z_{2}\right|$.

Ans:

- 5.2: If $x_{1}, x_{2} \in \mathbb{R}$ (are real numbers), define the meaning of $x_{1}>x_{2}$.


## Ans:

- 5.3: Explain the meaning of $z_{1}>z_{2}$.


## Ans:

- 5.4: (2 pts) What is the meaning of $\left|z_{1}+z_{2}\right|>3$ ?


## Ans:

- 5.5: (2 pts) If time were complex, how might the world be different?

Ans:

Problem \# 6: (1 pt) It is frequently necessary to consider a function $w(z)=u+v \jmath$ in terms of the real functions $u(x, y)$ and $v(x, y)$ (e.g. separate the real and imaginary parts). Similarly, we can consider the inverse $z(w)=x+y$, where $x(u, v)$ and $y(u, v)$ are real functions.

- 6.1: (1 pts) Find $u(x, y)$ and $v(x, y)$ for $w(z)=1 / z$.


## Ans:

Problem \# 7: (5 pts) Find $u(x, y)$ and $v(x, y)$ for $w(z)=c^{z}$ with complex constant $c \in \mathbb{C}$ for questions 7.1, 7.2, and 7.3:

$$
-7.1: c=e
$$

Ans:

$$
\text { - 7.2: } c=1 \text { (recall that } 1=e^{ \pm 22 \pi k} \text { for } k \in \mathbb{Z}
$$

## Ans:

$$
-7.3: c=\jmath . \text { Hint: } \jmath=e^{\jmath \pi / 2+\jmath 2 \pi k}, \quad k \in \mathbb{Z}
$$

## Ans:

- 7.4: (2 pts) What is $\jmath^{?}$ ?


## Ans:



Figure 2.2: This figure shows how to derive the Schwarz inequality, by finding the value of $\alpha=\alpha^{*}$ corresponding to min $[\boldsymbol{E}(\alpha)]$. It is identical to Fig. 3.5 on page 91.

## Schwarz inequality

Problem \# 8: The above figure shows three vectors for an arbitrary value of $\alpha \in \mathbb{R}$ and a specific value of $\alpha=\alpha^{*}$.
-8.1: Find the value of $\alpha \in \mathbb{R}$ such that the length (norm) of $\vec{E}$ (i.e., $\|\vec{E}\| \geq 0$ ) is minimum. Show your derivation, not the answer ( $\alpha=\alpha^{*}$ ).

## Ans:

- 8.2: Find the formula for $\left\|\boldsymbol{E}\left(\alpha^{*}\right)\right\|^{2} \geq 0$. Hint: Substitute $\alpha^{*}$ into Eq. 3.5.9 (p. 92) and show that this results in the Schwarz inequality

$$
|\vec{U} \cdot \vec{V}| \leq\|\vec{U}\|\|\mid \vec{V}\| .
$$

## Ans:

Problem \# 9: Geometry and scaler products

- 9.1: What is the geometrical meaning of the dot product of two vectors?

Ans:

- 9.2: Give the formula for the dot product of two vectors. Explain the meaning based on Fig. 3.4 (page 87).


## Ans:

-9.3: Write the formula for the dot product of two vectors $\vec{U} \cdot \vec{V}$ in $\mathbb{R}^{n}$ in polar form (e.g., assume the angle between the vectors is $\theta$ ).

## Ans:

-9.4: How is the Schwarz inequality related to the Pythagorean theorem?

## Ans:

-9.5: Starting from $\|\boldsymbol{U}+\boldsymbol{V}\|$, derive the triangle inequality

$$
\|\vec{U}+\vec{V}\| \leq\|\vec{U}\|+\|\vec{V}\| .
$$

## Ans:

-9.6: The triangle inequality $\|\vec{U}+\vec{V}\| \leq\|\vec{U}\|+\|\vec{V}\|$ is true for two and three dimensions: Does it hold for five-dimensional vectors?

Ans:

- 9.7: Show that the wedge product $\vec{U} \wedge \vec{V} \perp \vec{U} \cdot \vec{V}$.

Ans:

## Probability

Problem \# 10: Basic terminology of experiments

- 10.1: What is the mean of a trial, and what is the average over all trials?

Ans:

- 10.2: What is the expected value of a random variable $X$ ?

Ans:

- 10.3: What is the standard deviation about the mean?

Ans:

- 10.4: What is the definition of information of a random variable?

Ans:

- 10.5: How do you combine events? Hint: If the event is the fip of a biased coin, the events are $H=p, T=1-p$, so the event is $\{p, 1-p\}$. To solve the problem, you must find the probabilities of two independent events.
Ans:
- 10.6: What does the term independent mean in the context of question 10.5? Give an example.

Ans:

- 10.7: Define odds.

Ans:

## Chapter 3

## Differential equations

### 3.1 Problems DE-1

### 3.1.1 Topics of this homework:

Complex numbers and functions (ordering and algebra), complex power series, fundamental theorem of calculus (real and complex); Cauchy-Riemann conditions, multivalued functions (branch cuts and Riemann sheets)

### 3.1.2 Complex Power Series

Problem \# 1: In each case derive (e.g., using Taylor's formula) the power series of $w(s)$ about $s=0$ and give the RoC of your series. If the power series doesn't exist, state why! Hint: In some cases, you can derive the series by relating the function to another function for which you already know the power series at $s=0$.

$$
-1.1: 1 /(1-s)
$$

## Ans:

$-1.2: 1 /\left(1-s^{2}\right)$

## Ans:

- 1.3: $1 /\left(1+s^{2}\right)$.


## Ans:

## $-1.4: 1 / s$

## Ans:

## $-1.5: 1 /\left(1-|s|^{2}\right)$

## Ans:

Problem \# 2: Consider the function $w(s)=1 / s$

- 2.1: Expand this function as a power series about $s=1$. Hint: Let $1 / s=1 /(1-1+s)=$ $1 /(1-(1-s))$.
Ans:
- 2.2: What is the RoC?

Ans:

- 2.3: Expand $w(s)=1 / s$ as a power series in $s^{-1}=1 / s$ about $s^{-1}=1$.

Ans:

- 2.4: What is the RoC?


## Ans:

- 2.5: What is the residue of the pole?

Ans:

Problem \# 3: Consider the function $w(s)=1 /(2-s)$

- 3.1: Expand $w(s)$ as a power series in $s^{-1}=1 / s$. State the RoC as a condition on $\left|s^{-1}\right|$. Hint: Multiply top and bottom by $s^{-1}$.
Ans:
- 3.2: Find the inverse function $s(w)$. Where are the poles and zeros of $s(w)$, and where is it analytic?


## Ans:

Problem \# 4:Summing the series
The Taylor series of functions have more than one region of convergence.
-4.1: Given some function $f(x)$, if $a=0.1$, what is the value of

$$
f(a)=1+a+a^{2}+a^{3}+\cdots ?
$$

Show your work. Ans:
-4.2: Let $a=10$. What is the value of

$$
f(a)=1+a+a^{2}+a^{3}+\cdots ?
$$

Ans:

### 3.1.3 Cauchy-Riemann Equations

Problem \# 5: For this problem $\jmath=\sqrt{-1}, s=\sigma+\omega \jmath$, and $F(s)=u(\sigma, \omega)+\jmath v(\sigma, \omega)$. According to the fundamental theorem of complex calculus (FTCC), the integration of a complex analytic function is independent of the path. It follows that the derivative of $F(s)$ is defined as

$$
\begin{equation*}
\frac{d F}{d s}=\frac{d}{d s}[u(\sigma, \omega)+\jmath v(\sigma, \omega)] . \tag{DE-1.1}
\end{equation*}
$$

If the integral is independent of the path, then the derivative must also be independent of the direction:

$$
\begin{equation*}
\frac{d F}{d s}=\frac{\partial F}{\partial \sigma}=\frac{\partial F}{\partial \jmath \omega} . \tag{DE-1.2}
\end{equation*}
$$

The Cauchy-Riemann (CR) conditions

$$
\frac{\partial u(\sigma, \omega)}{\partial \sigma}=\frac{\partial v(\sigma, \omega)}{\partial \omega} \quad \text { and } \quad \frac{\partial u(\sigma, \omega)}{\partial \omega}=-\frac{\partial v(\sigma, \omega)}{\partial \sigma}
$$

may be used to show where Equation DE-1.2 holds.
-5.1: Assuming Equation DE-1.2 is true, use it to derive the CR equations.

## Ans:

-5.2: Merge the CR equations to show that $u$ and $v$ obey Laplace's equations

$$
\nabla^{2} u(\sigma, \omega)=0 \quad \text { and } \quad \nabla^{2} v(\sigma, \omega)=0
$$

Ans:

What can you conclude?
Ans:

Problem \# 6: Apply the CR equations to the following functions. State for which values of $s=\sigma+i \omega$ the $C R$ conditions do or do not hold (e.g., where the function $F(s)$ is or is not analytic). Hint: Review where CR-1 and CR-2 hold.

$$
-6.1: F(s)=e^{s}
$$

## Ans:

-6.2: $F(s)=1 / s$
Ans:

### 3.1.4 Branch cuts and Riemann sheets

Problem \# 7: Consider the function $w^{2}(z)=z$. This function can also be written as $w_{ \pm}(z)=\sqrt{z_{ \pm}}$. Assume $z=r e^{\phi_{\jmath}}$ and $w(z)=\rho e^{\theta_{\jmath}}=\sqrt{r} e^{\phi_{\jmath} / 2}$.

- 7.1: How many Riemann sheets do you need in the domain $(z)$ and the range $(w)$ to fully represent this function as single-valued?
Ans:
- 7.2: Indicate (e.g., using a sketch) how the sheet(s) in the domain map to the sheet(s) in the range.
Ans:
- 7.3: Use zviz.m to plot the positive and negative square roots $+\sqrt{z}$ and $-\sqrt{z}$. Describe what you see.
Ans:
- 7.4: Where does zviz.mplace the branch cut for this function?


## Ans:

- 7.5: Must the branch cut necessarily be in this location?


## Ans:

Problem \# 8: Consider the function $w(z)=\log (z)$. As in Problem 7, let $z=r e^{\phi_{\jmath}}$ and $w(z)=\rho e^{\theta_{\jmath}}$. - 8.1: Describe with a sketch and then discuss the branch cut $f o r f(z)$.

Ans:
-8.2: What is the inverse of the function $z(f)$ ? Does this function have a branch cut? If so, where is it?
Ans:

- 8.3: Using zviz.m, show that

$$
\begin{equation*}
\tan ^{-1}(z)=-\frac{\jmath}{2} \log \frac{\jmath-z}{\jmath+z} . \tag{DE-1.3}
\end{equation*}
$$

In Fig. 4.1 (p. 134) these two functions are shown to be identical.
Ans:
-8.4: Algebraically justify Eq. DE-1.3. Hint: Let $w(z)=\tan ^{-1}(z)$ and $z(w)=\tan w=\sin w / \cos w$; then solve for $e^{w j}$.
Ans:

### 3.1.5 A Cauer synthesis of any Brune impedance

Problem \# 9: One may synthesize a transmission line (ladder network) from a positive real impedance $Z(s)$ by using the continued fraction method. To obtain the series and shunt impedance values, we can use a residue expansion. Here we shall explore this method.

- 9.1: Starting from the Brune impedance $Z(s)=\frac{1}{s+1}$, find the impedance network as a ladder network.


## Ans:

- 9.2: Use a residue expansion in place of the CFA floor function (Sec. 2.4.4, p. 30) for polynomial expansions. Find the residue expansion of $H(s)=s^{2} /(s+1)$ and express it as a ladder network.
Ans:
- 9.3: Discuss how the series impedance $Z(s, x)$ and shunt admittance $Y(s, x)$ determine the wave velocity $\kappa(s, x)$ and the characteristic impedance $z_{o}(s, x)$ when (1) $Z(s)$ and $Y(s)$ are both independent of $x$; (2) $Y(s)$ is independent of $x, Z(s, x)$ depends on $x$; (3) $Z(s)$ is independent of $x$, $Y(s, x)$ depends on $x$; and (4) both $Y(s, x), Z(s, x)$ depend on $x$.


## Ans:

### 3.2 Problems DE-2

### 3.2.1 Topics of this homework:

Integration of complex functions, Cauchy's theorem, integral formula, residue theorem, power series, Riemann sheets and branch cuts, inverse Laplace transforms, Quadratic forms.

## Two fundamental theorems of calculus

## Fundamental Theorem of Calculus (Leibniz):

According to the fundamental theorem of (real) calculus (FTC),

$$
\begin{equation*}
f(x)=f(a)+\int_{a}^{x} F(\xi) d \xi \tag{DE-2.1}
\end{equation*}
$$

where $x, a, \xi, F, f \in \mathbb{R}$. This is an indefinite integral (since the upper limit is unspecified). It follows that

$$
\frac{d f(x)}{d x}=\frac{d}{d x} \int_{a}^{x} F(x) d x=F(x) .
$$

This justifies also calling the indefinite integral the antiderivative.
For a closed interval $[a, b]$, the FTC is

$$
\begin{equation*}
\int_{a}^{b} F(x) d x=f(b)-f(a), \tag{DE-2.2}
\end{equation*}
$$

thus the integral is independent of the path from $x=a$ to $x=b$.

## Fundamental Theorem of Complex Calculus:

According to the fundamental theorem of complex calculus (FTCC),

$$
\begin{equation*}
f(z)=f\left(z_{0}\right)+\int_{z_{0}}^{z} F(\zeta) d \zeta \tag{DE-2.3}
\end{equation*}
$$

where $z_{0}, z, \zeta, f, F \in \mathbb{C}$. It follows that

$$
\begin{equation*}
\frac{d f(z)}{d z}=\frac{d}{d z} \int_{z_{0}}^{z} F(\zeta) d \zeta=F(z) . \tag{DE-2.4}
\end{equation*}
$$

For a closed interval $\left[s, s_{o}\right.$ ], the FTCC is

$$
\begin{equation*}
\int_{s_{o}}^{s} F(\zeta) d \zeta=f(s)-f\left(s_{o}\right), \tag{DE-2.5}
\end{equation*}
$$

thus the integral is independent of the path from $x=a$ to $x=b$.

## Problem \# I

- 1.1: Consider Equation DE-2.1. What is the condition on $F(x)$ for which this formula is true?

Ans:

- 1.2: Consider Equation DE-2.3. What is the condition on $F(z)$ for which this formula is true?


## Ans:

- 1.3: Let $F(z)=\sum_{k=0}^{\infty} c_{k} z^{k}$.


## Ans:

- 1.4: Let

$$
F(z)=\frac{\sum_{k=0}^{\infty} c_{k} z^{k}}{z-\jmath}
$$

Ans:

Problem \# 2: In the following problems, solve the integral

$$
I=\int_{\mathcal{C}} F(z) d z
$$

for a given path $\mathcal{C} \in \mathbb{C}$.

- 2.1: Perform the following integrals $(z=x+i y \in \mathbb{C})$ :
$I=\int_{0}^{1+\jmath} z d z$
Ans:
-2.2: $I=\int_{0}^{1+j} z d z$, but this time make the path explicit: from 0 to 1 , with $y=0$, and then to $y=1$, with $x=1$.

Ans:

- 2.3: Discuss whether your results agree with Eq. DE-2.4?


## Ans:

Problem \# 3: Perform the following integrals on the closed path $C$, which we define to be the unit circle. You should substitute $z=e^{i \theta}$ and $d z=i e^{i \theta} d \theta$, and integrate from $\{-\pi, \pi\}$ to go once around the unit circle.

Discuss whether your results agree with Eq. DE-2.4?
$-3.1: \int_{C} z d z$

## Ans:

-3.2: $\int_{\mathcal{C}} \frac{1}{z} d z$

## Ans:

-3.3: $\int_{C} \frac{1}{z^{2}} d z$

## Ans:

-3.4: $I=\int_{C} \frac{1}{(z+2 \jmath)^{2}} d z$.
Recall that the path of integration is the unit circle, starting and ending at -1 .
Ans:

Problem \# 4: FTCC and integration in the complex plane
Let the function $F(z)=c^{z}$, where $c \in \mathbb{C}$ is given for each question. Hint: Can you apply the FTCC?

- 4.1: For the function $f(z)=c^{z}$, where $c \in \mathbb{C}$ is an arbitrary complex constant, use the Cauchy-Riemann $(C R)$ equations to show that $f(z)$ is analytic for all $z \in \mathbb{C}$.


## Ans:

- 4.2: Find the antiderivative of $F(z)$.


## Ans:

-4.3: $c=1 / e=1 / 2.7183, \ldots$ where $\mathcal{C}$ is $\zeta=0 \rightarrow i \rightarrow z$

## Ans:

- 4.4: $c=2$, where $C$ is $\zeta=0 \rightarrow(1+i) \rightarrow z$


## Ans:

-4.5: $c=i$, where the path $C$ is an inward spiral described by $z(t)=0.99^{t} e^{i 2 \pi t}$ for $t=0 \rightarrow t_{0} \rightarrow \infty$

## Ans:

-4.6: $c=e^{t-\tau_{0}}$, where $\tau_{0}>0$ is a real number and $\mathcal{C}$ is $z=(1-i \infty) \rightarrow(1+i \infty)$. Hint: Do you recognize this integral? If you do not, please do not spend a lot of time trying to solve it via the "brute force" method.

## Ans:

### 3.2.2 Cauchy's theorems CT-1, CT-2, CT-3

There are three basic definitions related to Cauchy's integral formula. They are all related and can greatly simplify integration in the complex plane. When a function depends on a complex variable, we use uppercase notation, consistent with the engineering literature for the Laplace transform.

Problem \# 5: Describe the relationships between the theorems:

$$
-5.1: C T-1 \text { and CT-2 }
$$

## Ans:

- 5.2: CT-1 and CT-3

Ans:

- 5.3: CT-2 and CT-3

Ans:

- 5.4: Consider the function with poles at $z= \pm j$,

$$
F(z)=\frac{1}{1+z^{2}}=\frac{1}{(z-j)(z+j)}
$$

Find the residue expansion.
Ans:

Problem \# 6: Apply Cauchy's theorems to solve the following integrals. State which theorem(s) you used and show your work.

- 6.1: $\oint_{\mathcal{C}} F(z) d z$, where $C$ is a circle centered at $z=0$ with a radius of $\frac{1}{2}$


## Ans:

- 6.2: $\oint_{C} F(z) d z$, where $C$ is a circle centered at $z=j$ with a radius of 1


## Ans:

- 6.3: $\oint_{\mathcal{C}} F(z) d z$, where $C$ is a circle centered at $z=0$ with a radius of 2


## Ans:

## Integration of analytic functions

Problem \# 7: In the following questions, you'll be asked to integrate $F(s)=u(\sigma, \omega)+i v(\sigma, \omega)$ around the contour $C$ for complex $s=\sigma+i \omega$,

$$
\begin{equation*}
\oint_{\mathcal{C}} F(s) d s \tag{DE-2.6}
\end{equation*}
$$

Follow the directions carefully for each question. When asked to state where the function is and is not analytic, you are not required to use the Cauchy-Riemann equations

$$
-7.1: F(s)=\sin (s)
$$

Ans:

- 7.2: Given function $F(s)=\frac{1}{s}$ State where the function is and is not analytic.

Ans:

- 7.3: Explicitly evaluate the integral when $\mathcal{C}$ is the unit circle, defined as $s=e^{i \theta}, 0 \leq \theta \leq 2 \pi$. Ans:
- 7.4: Evaluate the same integral using Cauchy's theorem and/or the residue theorem.


## Ans:

- 7.5: $F(s)=\frac{1}{s^{2}}$ State where the function is and is not analytic.


## Ans:

- 7.6: Explicitly evaluate the integral when $\mathcal{C}$ is the unit circle, defined as $s=e^{i \theta}, 0 \leq \theta \leq 2 \pi$. Ans:
- 7.7: What does your result imply about the residue of the second-order pole at $s=0$ ?


## Ans:

- 7.8: $F(s)=e^{s t}$ : State where the function is and is not analytic.


## Ans:

-7.9: Explicitly evaluate the integral when $\mathcal{C}$ is the square $(\sigma, \omega)=(1,1) \rightarrow(-1,1) \rightarrow(-1,-1) \rightarrow(1,-1)$

## Ans:

- 7.10: Evaluate the same integral using Cauchy's theorem and/or the residue theorem.


## Ans:

- 7.11: $F(s)=\frac{1}{s+2}$ : State where the function is and is not analytic.


## Ans:

- 7.12: Let $\mathcal{C}$ be the unit circle, defined as $s=e^{i \theta}, 0 \leq \theta \leq 2 \pi$. Evaluate the integral using Cauchy's theorem and/or the residue theorem.


## Ans:

- 7.13: Let $\mathcal{C}$ be a circle of radius 3, defined as $s=3 e^{i \theta}, 0 \leq \theta \leq 2 \pi$. Evaluate the integral using Cauchy's theorem and/or the residue theorem.


## Ans:

- 7.14: $F(s)=\frac{1}{2 \pi i} \frac{e^{s t}}{(s+4)}$ State where the function is and is not analytic.


## Ans:

- 7.15: Let $\mathcal{C}$ be a circle of radius 3, defined as $s=3 e^{i \theta}, 0 \leq \theta \leq 2 \pi$. Evaluate the integral using Cauchy's theorem and/or the residue theorem.


## Ans:

- 7.16: Let C contain the entire left half s plane. Evaluate the integral using Cauchy's theorem and/or the residue theorem. Do you recognize this integral?


## Ans:

- 7.17: $F(s)= \pm \frac{1}{\sqrt{s}}\left(\right.$ e.g., $\left.F^{2}=\frac{1}{s}\right)$ State where the function is and is not analytic.


## Ans:

- 7.18: This function is multivalued. How many Riemann sheets do you need in the domain (s) and the range (f) to fully represent this function? Indicate (e.g., using a sketch) how the sheet(s) in the domain map to the sheet(s) in the range.


## Ans:

- 7.19: Explicitly evaluate the integral $\int_{\mathcal{C}} \frac{1}{\sqrt{z}} d z$ when $\mathcal{C}$ is the unit circle, defined as $s=e^{i \theta}$, $0 \leq \theta \leq 2 \pi$. Is this contour closed? State why or why not.


## Ans:

- 7.20: Explicitly evaluate the integral $\int_{\mathcal{C}} \frac{1}{\sqrt{z}} d z$ when $\mathcal{C}$ is twice around the unit circle, defined as $s=e^{i \theta}, 0 \leq \theta \leq 4 \pi$. Is this contour closed? State why or why not. Hint: Note that $\sqrt{e^{i(\theta+2 \pi)}}=\sqrt{e^{i 2 \pi}} e^{i \theta}=e^{i \pi} \sqrt{e^{i}}$


## Ans:

-7.21: What does your result imply about the residue of the (twice-around $\frac{1}{2}$ order) pole at $s=0$ ?

## Ans:

- 7.22: Show that the residue is zero. Hint: Apply the definition of the residue.


## Ans:

### 3.2.3 Laplace transform applications

Problem \# 8: A two-port network application for the Laplace transform


Figure 3.1: This three-element electrical circuit is a system that acts to low-pass filter the signal voltage $V_{1}(\omega)$, to produce signal $V_{2}(\omega)$. It is convenient to define the dimensionless ratio $s / s_{c}=R C$ s in terms of a time constant $\tau=R C$ and cutoff frequency $s_{c}=1 / \tau$.

- 8.1: Find the $2 \times 2$ ABCD matrix representation of Fig. 8. Express the results in terms of the dimensionless ratio $s / s_{c}$, where $s_{c}=1 / \tau$ is the cutoff frequency and $\tau=R C$ is the time constant. Ans:
-8.2: Find the eigenvalues of the $2 \times 2$ ABCD matrix. Hint: See Appendix B.3.
Ans:
-8.3: Assuming that $I_{2}=0$, find the transfer function $H(s) \equiv V_{2} / V_{1}$. From the results of the ABCD matrix you determined in questions ?? and ??, show that

$$
\begin{equation*}
H(s)=\frac{s_{c}}{s+s_{c}} . \tag{DE-2.7}
\end{equation*}
$$

## Ans:

- 8.4: The transfer function $H(s)$ has one pole. Where is the pole and residue?


## Ans:

- 8.5: Find $h(t)$, the inverse Laplace transform of $H(s)$.


## Ans:

- 8.6: Assuming that $V_{2}=0$, find $Y_{12}(s) \equiv I_{2} / V_{1}$.


## Ans:

- 8.7: Find the input impedance to the right-hand side of the system $Z_{22}(s) \equiv V_{2} / I_{2}$ for two cases: (1) $I_{1}=0$ and (2) $V_{1}=0$.


## Ans:

- 8.8: Compute the determinant of the $A B C D$ matrix. Hint: It is always $\pm 1$.


## Ans:

- 8.9: Given the result of the previous problem Eq. ??, compute the derivative of $H(s)=\left.\frac{V_{2}}{V_{1}}\right|_{I_{2}=0}$.

Ans:

### 3.2.4 Computer exercises with Matlab/Octave

Problem \# 9: With the help of a computer
Now we look at a few important concepts using Matlab/Octave's syms commands or Wolfram Alpha's symbolic math toolbox. ${ }^{1}$

For example, to find the Taylor series expansion about $s=0$ of

$$
F(s)=-\log (1-s)
$$

we first consider the derivative and its Taylor series (about $s=0$ )

$$
F^{\prime}(s)=\frac{1}{1-s}=\sum_{n=0}^{\infty} s^{n} .
$$

Then, we integrate this series term by term:

$$
F(s)=-\log (1-s)=\int^{s} F^{\prime}(s) d s=\sum_{n=0}^{\infty} \frac{s^{n}}{n}
$$

Alternatively we can use Matlab/Octave commands:

```
syms s
taylor(-log(1-s),'order', 7)
```

- 9.1: Use Octave's taylor (-log (1-s)) to the seventh order, as in the example above. Try the above Matlab/Octave commands. Give the first seven terms of the Taylor series (confirm that Matlab/Octave agrees with the formula derived above).


## Ans:

-9.2: What is the inverse Laplace transform of this series? Consider the series term by term.

## Ans:

- 9.3: The function $1 / \sqrt{z}$ has a branch point at $z=0$; thus it is singular there. Can you apply Cauchy's integral theorem when integrating around the unit circle?


## Ans:

[^5]- 9.4: This Matlab/Octave code computes $\int_{0}^{4 \pi} \frac{d z}{\sqrt{z}}$ using Matlab's/Octave's symbolic analysis package. Run the following script:

```
syms z
I=int(1/sqrt(z))
J = int(1/sqrt(z), exp(-j*pi), exp(j*pi))
eval(J)
```

What answers do you get for $I$ and $J$ ?
Ans:
-9.5: Modify this code to integrate $f(z)=1 / z^{2}$ once around the unit circle. What answers do you get for $I$ and $J$ ?
Ans:

- 9.6: Bessel functions can describe waves in a cylindrical geometry.

The Bessel function has a Laplace transform with a branch cut

$$
J_{0}(t) u(t) \leftrightarrow \frac{1}{\sqrt{1+s^{2}}} .
$$

Draw a hand sketch showing the nature of the branch cut. Hint: Use zviz. Ans:

Problem \# 10: Matlab/Octave exercises

- 10.1: Try the following Matlab/Octave commands, and then comment on your findings.

```
%Take the inverse LT of 1/sqrt(1+s^2)
syms s
I=ilaplace(1/(sqrt((1+s^2))));
disp(I)
```


## Ans:

\%Find the Taylor series of the LT
$\mathrm{T}=$ taylor(1/sqrt(1+s^2),10); disp(T);

## Ans:

```
%Verify this
```

syms t
J=laplace (bessel $j(0, t))$;
disp(J);

## Ans:

\%plot the Bessel function $\mathrm{t}=0: 0.1: 10 * \mathrm{pi}$;
b=besselj(0,t); plot(t/pi,b); grid on;
Ans:

- 10.2: When did Friedrich Bessel live?


## Ans:

- 10.3: What did he use Bessel functions for?

Ans:

Problem \# 11: Colorized plots of analytic functions. Use zviz for each of the following.

- 11.1: Describe the plot generated by zviz $S=z$.


## Ans:

- 11.2: Describe the plot generated by zviz 1./sqrt( $\left.1+S^{2}{ }^{2}\right)$.


## Ans:

- 11.3: Describe the plot generated by zviz 1./sqrt(1-S. ${ }^{2}$ ). Is this function a Brune impedance (i.e., does this function obey


## Ans:

## - 11.4: zviz 1./(1+sqrt(S))

## Ans:

### 3.2.5 Inverse of Riemann $\zeta(s)$ function

Problem \# 12: Inverse zeta function (This problem is for extra credit).

- 12.1: Find the $\mathcal{L T}{ }^{-1}$ of one factor of the Riemann zeta function $\zeta_{p}(s)$, where $\zeta_{p}(s) \leftrightarrow z_{p}(t)$. Describe your results in words. Hint: See Eq. ??. Hint: Consider the geometric series representation

$$
\begin{equation*}
\zeta_{p}(s)=\frac{1}{1-e^{-s T_{p}}}=\sum_{k=0}^{\infty} e^{-s k T_{p}}, \tag{DE-2.8}
\end{equation*}
$$

for which you can look up the $L \mathcal{T}^{-1}$ of each term.
Ans:

Problem \# 13: Inverse transform of products:
The time-domain version of Eq. DE-2.7 may be written as the convolution of all the $z_{k}(t)$ factors:

$$
\begin{equation*}
z(t) \equiv z_{2}(t) \star z_{3}(t) \star z_{5}(t) \star z_{7}(t) \star \cdots \star z_{p}(t) \star \cdots, \tag{DE-2.9}
\end{equation*}
$$

where $\star$ represents time convolution.

Explain what this means in physical terms. Start with two terms (e.g., $z_{1}(t) \star z_{2}$ ). Hint: The input admittance of this cascade may be interpreted as the analytic continuation of $\zeta(s)$ by defining a cascade of eigenfunctions with eigenvalues derived from the primes. For a discussion of this idea see Sec. 3.2.3 and C.1.1.

Ans:

Physical interpretation: Such functions may be generated in the time domain, as shown in Fig. ??, using a feedback delay of $T_{p}$ seconds, described by the two equations in the Fig. ?? with a unity feedback gain $\alpha=-1$. Taking the Laplace transform of the system equation, we see that the transfer function between the state variable $q(t)$ and the input $x(t)$ is given by $\zeta_{p}(s)$, which is an all-pole function, since

$$
\begin{equation*}
Q(s)=e^{-s T_{n}} Q(s)+V(s), \text { or } \zeta_{p}(s) \equiv \frac{Q(s)}{V(s)}=\frac{1}{1-e^{-s T_{p}}} \tag{DE-2.10}
\end{equation*}
$$

Closing the feed-forward path gives a second transfer function $Y(s)=I(s) / V(s)$-namely,

$$
\begin{equation*}
Y(s) \equiv \frac{I(s)}{V(s)}=\frac{1-e^{-s T_{p}}}{1+e^{-s T_{p}}} \tag{DE-2.11}
\end{equation*}
$$

If we take $i(t)$ as the current and $v(t)$ as the voltage at the input to the transmission line, then $y_{p}(t) \leftrightarrow \zeta_{p}(s)$ represents the input impedance at the input to the line. The poles and zeros of the impedance interleave along the $j \omega$ axis. By a slight modification, $\zeta_{p}(s)$ may alternatively be written as

$$
\begin{equation*}
Y_{p}(s)=\frac{e^{s T_{p} / 2}+e^{-s T_{p} / 2}}{e^{s T_{p} / 2}-e^{-s T_{p} / 2}}=j \tan \left(s T_{p} / 2\right) \tag{DE-2.12}
\end{equation*}
$$

Every impedance $Z(s)$ has a corresponding reflectance function given by a Möbius transformation, which may be read off of Eq. DE-2.11 as

$$
\begin{equation*}
\Gamma(s) \equiv \frac{1+Z(s)}{1-Z(s)}=e^{-s T_{p}} \tag{DE-2.13}
\end{equation*}
$$

since impedance is also related to the round-trip delay $T_{p}$ on the line. The inverse Laplace transform of $\Gamma(s)$ is the round-trip delay $T_{p}$ on the line

$$
\begin{equation*}
\gamma(t)=\delta\left(t-T_{p}\right) \leftrightarrow e^{-s T_{p}} \tag{DE-2.14}
\end{equation*}
$$

Working in the time domain provides a key insight, as it allows us to parse out the best analytic continuation of the infinity of possible continuations that are not obvious in the frequency domain (See Sec. 3.2.3). Transforming to the time domain is a form of analytic continuation of $\zeta(s)$ that depends on the assumption that $Z^{\text {eta }}(t) \leftrightarrow \zeta(s)$ is one-sided in time (causal).

### 3.2.6 Quadratic forms

A matrix that has positive eigenvalues is said to be positive-definite. The eigenvalues are real if the matrix is symmetric, so this is a necessary condition for the matrix to be positive-definite. This condition is related to conservation of energy, since the power is the voltage times the current. Given an impedance matrix

$$
\mathbf{V}=z \mathbf{I}
$$

the power $\mathcal{P}$ is

$$
\mathcal{P}=\mathbf{I} \cdot \mathbf{V}=\mathbf{I} \cdot Z \mathbf{I}
$$

which must be positive-definite for the system to obey conservation of energy.
Problem \# 14: In this problem, consider the $2 \times 2$ impedance matrix

$$
Z=\left[\begin{array}{ll}
2 & 1 \\
1 & 4
\end{array}\right]
$$

- 14.1: Solve for the power $\mathcal{P}\left(i_{1}, i_{2}\right)$ by multiplying out this matrix equation (which is a quadratic form):

$$
\mathscr{P}\left(i_{1}, i_{2}\right)=\mathbf{I}^{T}\left[\begin{array}{ll}
2 & 1 \\
1 & 4
\end{array}\right] \mathbf{I} .
$$

## Ans:

- 14.2: Is the impedance matrix positive-definite? Show your work by finding the eigenvalues of the matrix $\mathbf{Z}$.
Ans:
- 14.3: Should an impedance matrix always be positive-definite? Explain.

Ans:

### 3.3 Problems DE-3

### 3.3.1 Topics of this homework: Brune impedance

lattice transmission line analysis

### 3.3.2 Brune Impedance

Problem \# 1: Residue form
A Brune impedance is defined as the ratio of the force $F(s)$ to the flow $V(s)$ and may be expressed in residue form as

$$
\begin{equation*}
Z(s)=c_{0}+\sum_{k=1}^{K} \frac{c_{k}}{s-s_{k}}=\frac{N(s)}{D(s)} \tag{DE-3.1}
\end{equation*}
$$

with

$$
D(s)=\prod_{k=1}^{K}\left(s-s_{k}\right) \quad \text { and } \quad c_{k}=\lim _{s \rightarrow s_{k}}\left(s-s_{k}\right) D(s)=\prod_{n^{\prime}=1}^{K-1}\left(s-s_{n}\right)
$$

The prime on the index $n^{\prime}$ means that $n=k$ is not included in the product.

- 1.1: Find the Laplace transform (LT) of a (1) spring, (2) dashpot, and (3) mass.

Express these in terms of the force $F(s)$ and the velocity $V(s)$, along with the electrical equivalent impedance: (1) Hooke's law $f(t)=K x(t)$, (2) dashpot resistance $f(t)=R v(t)$, and (3) Newton's law for mass $f(t)=M d v(t) / d t$. Ans:

- 1.2: Take the Laplace transform $(\mathcal{L T})$ of Eq. DE-3.2 and find the total impedance $Z(s)$ of the mechanical circuit.

$$
\begin{equation*}
M \frac{d^{2}}{d t^{2}} x(t)+R \frac{d}{d t} x(t)+K x(t)=f(t) \leftrightarrow\left(M s^{2}+R s+K\right) X(s)=F(s) \tag{DE-3.2}
\end{equation*}
$$

Ans:

- 1.4: Assume that $M=R=K=1$ and find the residue form of the admittance $Y(s)=1 / Z(s)$ (see Eq. DE-3.1) in terms of the roots $s_{ \pm}$. Hint: Check your answer with Octave's/Matlab's res idue command.


## Ans:

- 1.5: By applying Eq. 4.5 .3 (page 151), find the inverse Laplace transform $\left(\mathcal{L} \mathcal{T}^{-1}\right)$. Use the residue form of the expression that you derived in question 1.4.
Ans:


### 3.3.3 Transmission-line analysis

Problem \# 2: Train-mission-line We wish to model the dynamics of a freight train that has $N$ such cars and study the velocity transfer function under various load conditions.

As shown in Fig. 4.8.2, the train model consists of masses connected by springs.


Figure 3.2: Depiction of a train consisting of cars treated as masses $M$ and linkages treated as springs of stiffness $K$ or compliance $C=1 / K$. Below it is the electrical equivalent circuit for comparison. The masses are modeled as inductors and the springs as capacitors to ground. The velocity is analogous to a current and the force $f_{n}(t)$ to the voltage $\phi_{n}(t)$. The length of each cell is $\Delta[\mathrm{m}]$. The train may be accurately modeled as a transmission line (TL), since the equivalent electrical circuit is a lumped model of a TL. This method, called a Cauer synthesis, is based on the ABCD transmission line method of Sec. 3.8 (p. 108).

Use the ABCD method (see the discussion in Appendix B.3, p. 228) to find the matrix representation of the system of Fig. 4.8.2. Define the force on the $n$th train car $f_{n}(t) \leftrightarrow F_{n}(\omega)$ and the velocity $v_{n}(t) \leftrightarrow V_{n}(\omega)$.

Break the model into cells consisting of three elements: a series inductor representing half the mass ( $M / 2$ ), a shunt capacitor representing the spring $(C=1 / K)$, and another series inductor representing half the mass ( $L=M / 2$ ), transforming the model into a cascade of symmetric $(\mathcal{A}=\mathcal{D})$ identical cell matrices $\mathcal{T}(s)$.

- 2.1: Find the elements of the $A B C D$ matrix $\mathcal{T}$ for the single cell that relate the input node 1 to output node 2

$$
\left[\begin{array}{c}
F  \tag{DE-3.3}\\
V
\end{array}\right]_{1}=\mathcal{T}\left[\begin{array}{c}
F(\omega) \\
-V(\omega)
\end{array}\right]_{2}
$$

Ans:

- 2.2: Express each element of $\mathcal{T}(s)$ in terms of the complex Nyquist ratio $s / s_{c}<1(s=2 \pi j f$, $s_{c}=2 \pi j f_{c}$ ). The Nyquist wavelength sampling condition is $\lambda_{c}>2 \Delta$. It says the critical wavelength $\lambda_{c}>2 \Delta$. Namely it is defined in terms the minimum number of cells $2 \Delta$, per minimum wavelength $\lambda_{c}$. The Nyquist wavelength sampling theorem says that there are at least two cars per wavelength.

Proof: From the figure, the distance between cars $\Delta=c_{o} T_{o}[\mathrm{~m}]$, where

$$
c_{o}=\frac{1}{\sqrt{M C}} \quad[\mathrm{~m} / \mathrm{s}]
$$

The cutoff frequency obeys $f_{c} \lambda_{c}=c_{o}$. The Nyquist critical wavelength is $\lambda_{c}=c_{o} / f_{c}>2 \Delta$. Therefore the Nyquist sampling condition is

$$
\begin{equation*}
f<f_{c} \equiv \frac{c_{o}}{\lambda_{c}}=\frac{c_{o}}{2 \Delta}=\frac{1}{2 \Delta \sqrt{M C}} \quad[\mathrm{rad} / \mathrm{sec}] \tag{DE-3.4}
\end{equation*}
$$

Finally, $s_{c}=\jmath 2 \pi f_{c}$.
Ans:

- 2.3: Use the property of the Nyquist sampling frequency $\omega<\omega_{c}$ (Eq. DE-3.4) to remove higher order powers of frequency

$$
\begin{equation*}
1+\left(\frac{s}{s_{c}}\right)^{z^{0}} \approx 1 \tag{DE-3.5}
\end{equation*}
$$

to determine a band-limited approximation of $\mathcal{T}(s)$.
Ans:

Problem \# 3: Now consider the cascade of $N$ such $\mathcal{T}(s)$ matrices and perform an eigenanalysis.

- 3.1: Find the eigenvalues and eigenvectors of $\mathcal{T}(s)$ as functions of $s / s_{c}$.

Ans:

Problem \# 4: Find the velocity transferfunction $H_{12}(s)=V_{2} /\left.V_{1}\right|_{F_{2}=0}$.

- 4.1: Assuming that $N=2$ and $F_{2}=0$ (two half-mass problem), find the transfer function $H(s) \equiv V_{2} / V_{1}$. From the results of the $\mathcal{T}$ matrix, find

$$
H_{21}(s)=\left.\frac{V_{2}}{V_{1}}\right|_{F_{2}=0}
$$

Express $H_{12}$ in terms of a residue expansion. Ans:

- 4.2: Find $h_{21}(t) \leftrightarrow H_{21}(s)$.

Ans:
-4.3: What is the input impedance $Z_{2}=F_{2} / V_{2}$, assuming $F_{3}=-r_{0} V_{3}$ ?
Ans:

- 4.4: Simplify the expression for $Z_{2}$ as follows:

1. Assuming the characteristic impedance $r_{0}=\sqrt{M / C}$,
2. terminate the system in $r_{0}: F_{2}=-r_{0} V_{2}$ (i.e., $-V_{2}$ cancels).
3. Assume higher-order frequency terms are less than $1\left(\left|s / s_{c}\right|<1\right)$.
4. Let the number of cells $N \rightarrow \infty$. Thus $\left|s / s_{c}\right|^{N}=0$.

When a transmission line is terminated in its characteristic impedance $r_{0}$, the input impedance $Z_{1}(s)=r_{0}$. Thus, when we simplify the expression for $\mathcal{T}(s)$, it should be equal to $r_{0}$. Show that this is true for this setup.

Ans:

- 4.5: State the ABCD matrix relationship between the first and Nth nodes in terms of the cell matrix. Write out the transfer function for one cell, $H_{21}$.
Ans:
-4.6: What is the velocity transfer function $H_{N 1}=\frac{V_{N}}{V_{1}}$ ?
Ans:


[^0]:    ${ }^{1}$ This problem came from an exam problem for Math 213, Fall 2016.

[^1]:    ${ }^{2}$ This notation is highly nonstandard due to the nonlinear operations. The matrix elements are derived from the vector rather than multiplying them. These calculation may be done with the help of Matlab/Octave

[^2]:    ${ }^{1}$ The geometric series is not defined as the function $1 /(1-x)$, it is defined as the series $1+x+x^{2}+x^{3}+\cdots$, such that the ratio of consecutive terms is $x$.

[^3]:    ${ }^{2}$ Taken from Rotman (1996, p. 50)

[^4]:    ${ }^{3}$ Greenberg p. 694, Eq. 8.

[^5]:    ${ }^{1}$ https://www.wolframalpha.com/

