

Experiments in Physics

PHYSICS 1292
GENERAL PHYSICS II LAB

COLUMBIA UNIVERSITY



DEPARTMENT OF PHYSICS

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Introduction 2-0

General Instructions

1 Purpose of the Laboratory

The laboratory experiments described in this manual are an important part of your physics course. Most of the experiments are designed to illustrate important concepts described in the lectures. Whenever possible, the material will have been discussed in lecture before you come to the laboratory.

The sections headed Applications and Lab Preparation Examples, which are included in some of the manual sections, are *not* required reading unless your laboratory instructor specifically assigns some part. The Applications are intended to be motivational and so should indicate the importance of the laboratory material in medical and other applications. The Lab Preparation Examples are designed to help you prepare for the lab; you will not be required to answer all these questions (though you should be able to answer any of them by the end of the lab). The individual laboratory instructors may require you to prepare answers to a subset of these problems.

2 Preparation for the Laboratory

In order to keep the total time spent on laboratory work within reasonable bounds, the write-up for each experiment will be completed at the end of the lab and handed in *before the end of each laboratory period*. Therefore, it is imperative that you spend sufficient time preparing for the experiment *before* coming to laboratory. You should take advantage of the opportunity that the experiments are set up in the Lab Library (Room 506) and that TAs there are willing to discuss the procedure with you.

At each laboratory session, the instructor will take a few minutes at the beginning to go over the experiment to describe the equipment to be used and to outline the important issues. This does not substitute for careful preparation beforehand! You are expected to have studied the manual and appropriate references at home so that you are prepared when you arrive to perform the experiment. The instructor will be available primarily to answer questions, aid you in the use of the equipment, discuss the physics behind the experiment, and guide you in completing your analysis and write-up. Your instructor will describe his/her policy regarding expectations during the first lab meeting.

Some experiments and write-ups may be completed in less than the three-hour laboratory period, but under no circumstances will you be permitted to stay in the lab

after the end of the period or to take your report home to complete it. If it appears that you will be unable to complete all parts of the experiment, the instructor will arrange with you to limit the experimental work so that you have enough time to write the report during the lab period.

Note: Laboratory equipment must be handled with care and each laboratory bench must be returned to a neat and orderly state before you leave the laboratory. In particular, you must turn off all sources of electricity, water, and gas.

3 Bring to Each Laboratory Session

- A pocket calculator (with basic arithmetic and trigonometric operations).
- A pad of 8.5×11 inch graph paper and a sharp pencil. (You will write your reports on this paper, including your graphs. Covers and staplers will be provided in the laboratory.)
- A ruler (at least 10 cm long).

4 Graph Plotting

Frequently, a graph is the clearest way to represent the relationship between the quantities of interest. There are a number of conventions, which we include below.

- A graph indicates a relation between two quantities, x and y , when other variables or parameters have fixed values. Before plotting points on a graph, it may be useful to arrange the corresponding values of x and y in a table.
- Choose a convenient scale for each axis so that the plotted points will occupy a substantial part of the graph paper, but do not choose a scale which is difficult to plot and read, such as 3 or $3/4$ units to a square. Graphs should usually be at least half a page in size.
- Label each axis to identify the variable being plotted and the units being used. Mark prominent divisions on each axis with appropriate numbers.
- Identify plotted *points* with appropriate symbols, such as crosses, and when necessary draw vertical or horizontal *bars* through the points to indicate the range of uncertainty involved in these points.
- Often there will be a theory concerning the relationship of the two plotted variables. A linear relationship can be demonstrated if the data points fall along a

single straight line. There are mathematical techniques for determining which straight line best fits the data, but for the purposes of this lab it will be sufficient if you simply make a rough estimate visually. *The straight line should be drawn as near the mean of the all various points as is optimal.* That is, the line need not precisely pass through the first and last points. Instead, each point should be considered as accurate as any other point (unless there are experimental reasons why some points are less accurate than others). The line should be drawn with about as many points above it as below it, and with the ‘aboves’ and ‘below’ distributed at random along the line. (For example, not all points should be above the line at one end and below at the other end).

5 Questions or Complaints

If you have a difficulty, you should attempt to work it through with your laboratory instructor. If you cannot resolve it, you may discuss such issues with:

- one of the laboratory Preceptors in Pupin Room 729;
- the Undergraduate Assistant in the Departmental Office – Pupin Room 704;
- the instructor in the lecture course, or the Director of Undergraduate Studies;
- your undergraduate advisor.

As a general rule, it is a good idea to work downward through this list, though some issues may be more appropriate for one person than another.

Experiment 2-1

Electric Fields

1 Introduction

One of the fundamental concepts used to describe electric phenomena is that of the electric field. (Later we will also deal with magnetic fields.) We explore the field concept experimentally by measuring lines of equal electric potential on carbon paper and then constructing the electric field lines that connect them.

2 Theory

2.1 Fields

A field is a function in which a numerical value associated with a physical quantity is assigned to every point in space. For example, each point in the United States may be assigned a temperature (that we measure with a thermometer at that position), as shown in figure 1. Thus, temperature as a function of location is called the “temperature field”. The temperature field is an example of a scalar field, since the value assigned to each location is a scalar.

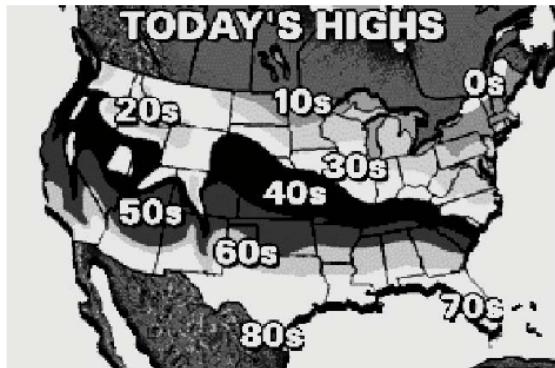


Figure 1: Temperature of Different Places in the United States

The electric field (as well as the magnetic field) is a vector field, because a vector (rather than a scalar) is associated with each position. At each point in space, the magnitude and direction of the electric field are equal to the magnitude and direction of the electric force that would act on a unit test charge at this position.

A line that is formed by connecting electric field vectors and follows the direction

of the field is called a field line. A field is considered homogeneous if the field lines are parallel. The potential associated with such a field increases linearly with the distance traveled along the field lines.

2.2 Equipotential Contours for Gravity

A topographical map is a map with lines that indicate contours of the same height. On such a map, a mountain may look like figure 2

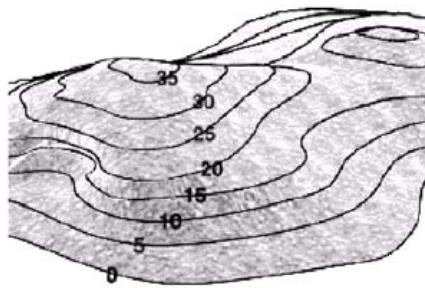


Figure 2: A Topographical Map

These contours of the same height (h) are equi-height lines or, if you think in terms of potential energy ($U = mgh$), they are equipotential lines. That is, every point on the contour has the same value of potential energy. So as you move along a contour, you do not change height, and therefore you neither gain nor lose potential energy. Only as you step up or down the hill do you change the energy.

There are a few general properties of equipotential contours:

- They never intersect. (A single point cannot be at two different heights at the same time, and therefore it cannot be on two different contours.)
- They close on themselves. (A line corresponding to constant height cannot just end.)
- They are smooth curves, as long as the topography has no sharp discontinuities (like a cliff).

2.3 Analogy between Electric and Gravitational Potentials

Electric potential is analogous to gravitational potential energy. Of course, gravitational potential energy arises from any object (with mass), whereas the electric potential arises only from charged objects.

As in mechanics, the absolute value of potential is not important. The differences in potential are the quantities that have physical meaning. In a real problem, it is usually best to choose a reference point, define it as having zero potential, and refer to potentials at all other points relative to that point. Similarly, topographic height is usually reported relative to the reference at sea level.

2.4 Field Lines

Let's exploit our analogy of mountains and valleys a bit further for electric fields. Electric field lines indicate the electric force on a test charge (magnitude and direction). Gravitational field lines tell us about the gravitational force on a test charge (magnitude and direction) – the direction it tends to roll down the hill and the acceleration it experiences during its journey. So if we carefully track a marble as we roll it down the hill in small increments, we can obtain a “field-line”.

What is the characteristic of a field line? Think about placing a marble on an inclined hillside. Which way will it roll if you release it? It will always roll in the direction of steepest descent, and this direction is always perpendicular to the contours that indicate equal height (which we have been referring to as equipotential lines).

So it is for the electric case: if we know the equipotential lines, we can draw the field lines such that each one is perpendicular to each equipotential line. Figure 3 and 4 both show equipotential lines (dashed lines) and the corresponding field lines (solid lines).

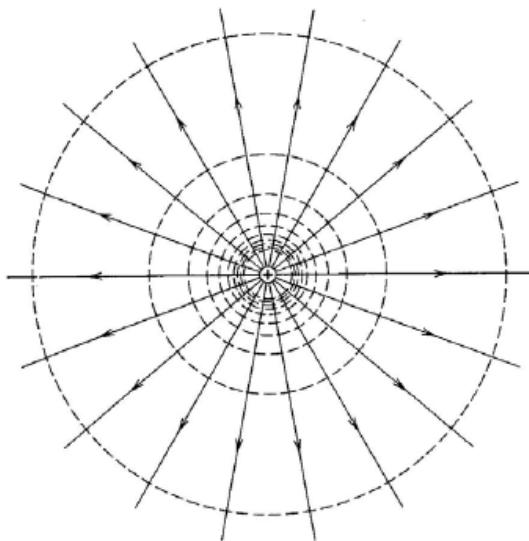


Figure 3: Equipotential Lines and Field Lines of a Single Charge

Keep in mind that the field lines (perpendicular to equipotential lines) give the biggest change in height over the shortest distance!

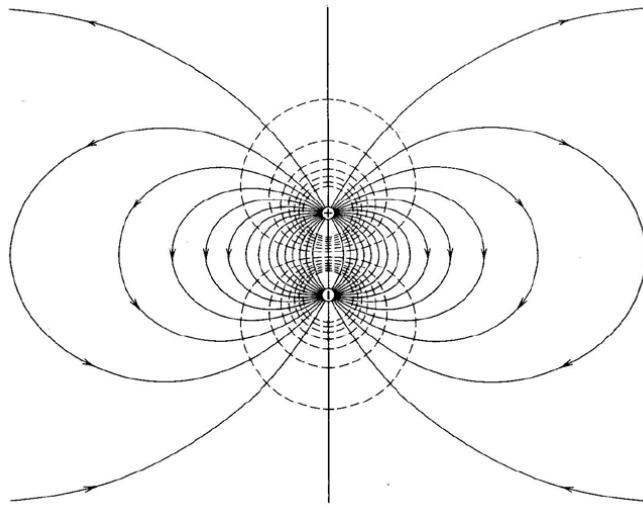


Figure 4: Equipotential Lines and Field Lines of a Pair of Opposite Charges

There are also a few rules for field lines:

- Field lines always begin and end on charges. (They often terminate on material surfaces, but that is because there are charges on those surfaces.)
- They never intersect each other (except at electric charges, where they also terminate).
- They always intersect equipotential lines perpendicularly!
- They are usually smoothly continuous (except if they terminate).

2.5 Metal Surfaces

It turns out that for electrical phenomena, metals are equipotentials. So, by using metals, we can impose some unusually shaped equipotential lines and determine the corresponding field lines.

Metals are always equipotentials because:

1. they are excellent conductors of electric charge, and
2. they have an abundant supply of freely moving negative charges (electrons).

Based on these two properties, we can understand that, if a field line were to penetrate a metal surface, these electrons would feel electric forces to move them through the material in the direction opposite to that of the electric field (since they are negative) until they would be forced to stop upon encountering the metal surface on the other side. Therefore, the free electrons ultimately distribute themselves along the surface of the metal so that there is a net negative charge on one side of the conductor and a net positive charge on the other (from where the electrons have fled). This realignment of the free electrons creates a field that precisely cancels the imposed field resulting in no net field, and the remaining free charges inside the conductor then feel no force and don't move! Hence, field lines imposed from outside always end on metal surfaces (which they intersect perpendicularly), and the surface (with the entire interior) is an equipotential.

2.6 Electric Shielding Theorem

As just discussed, whenever a piece of metal is placed in an electric field, the entire metal will remain an equipotential. That is, every point in the metal will be at the same potential as every other point. No field lines penetrate through metals.

What happens if we take an enclosed container of metal of arbitrary shape, say a tin can, and put it into an electric field? Since no field gets through the metal, there must be no electric field inside the can. This means that every point inside the container, as well as all points on the container surface must be at the same potential. (If they were not, then there would be differences in the potential and therefore an electric field.)

Lets return to our analogy of gravity equipotentials in hills and valleys, with field lines in the direction a marble would roll down the hill. Consider how a frozen lake would look on such a contour map. The surface of the lake has the same gravitational potential energy at all points, therefore if you place a marble on the frozen surface of the lake, the marble will not roll anywhere. In other words, there is no component of the gravitational field along the surface of the lake to push the marble. Similarly, there is no electric field on the inside of a closed metal container to push the charges around, and the entire interior is therefore a constant equipotential surface.

2.7 Remarks for Experts

1. The gravitational analogy operates in only two dimensions: the horizontal coordinates describing the surface of the lake. When discussing electrical fields we should, in principle, take into account all three dimensions. The analogy of the lake surface for the electrical case is the entire volume (three dimensions) of the interior of the can. For this experiment, we “cheat” a little by looking at electric field within a two dimensional world of slightly conducting paper. But the field

and potential arrangements are as one expects for electrostatic phenomena in a two dimensional system.

2. To shield electric field, we actually don't need a metallic surface that is completely closed; a closed cage of wire mesh (Faraday cage) is sufficient. So a sedan will work like a Faraday cage, but a convertible will not since it is not closed at the top.

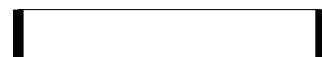
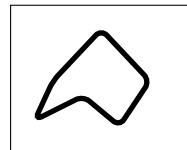
3 Experiments

In this experiment, we use pieces of slightly conductive paper, electrodes, and metallic paint to determine the equipotential and field configurations. Before we actually start our measurements we will have to draw the shapes on the conductive paper using the metal paint. The metal paint will take a few minutes to dry.

3.1 Preparation of the Conductive Paper

Take four pieces of conductive paper – three full-sized and one thin strip and paint on them with the metallic paint according to the following instructions. Each shape drawn by the metal paint is an equipotential curve. By connecting the pins to areas covered with paint, we can determine the field configurations between different equipotential shapes.

1. Leave the first piece of paper blank.
2. On the second piece of paper paint two parallel plates. Use a piece of cardboard provided to get sharp edges with the paint. Don't use the rulers!
3. On the third piece of paper paint a closed loop of any shape that does not enclose an electrode.
4. Take the thin strip and paint a line on each short end of the strip.



Hints:

1. You get the best results if you stay at least 1-2 cm away from the edges of the paper. Also make sure that the pins are always in contact with the paper. (Otherwise the experiment does not work!) Be careful not to make the holes too big or the pins will no longer be in contact with the paper.

2. Give the metal paint enough time to dry. First paint all the configurations you will use on the conductive paper. While they dry, get started on the two point-charge arrangement!

3.2 Setup of the Voltage Source

- “OUTPUT A” of the power supply will be used in this experiment to provide a constant potential difference of 10 volts.
- Turn the “A VOLTAGE”, the “A CURRENT”, the “B VOLTAGE”, and the “B CURRENT” control knobs located on the right side of the front panel counter-clockwise to the “MIN” position.
- Set the “A/B OUTPUT” switch located on the upper right of the front panel to the “INDEPENDENT” setting, and set the “A/B METER” switch located between the two meters on the front panel to the “A” setting.
- Turn the power on by using the button located on the left of the front panel.
- Slowly turn the “A VOLTAGE” output knob until the voltage reads 10 volts. After doing this, do not make any further changes on the power supply for the remainder of the experiment. (The current meter will stay at a zero reading.)
- Connect two cables to the “A OUTPUT (0 ~ 20 V 0.5 A)” of the power supply (red and black poles) and then connect these cables to two of the yellow-metallic pins.

NOTE: To use the digital multimeter as an appropriate voltmeter for this experiment: turn the multimeter power on, select “DC” (button in out position), select “VOLTS” (push button in), set range to “2 VOLTS”, and connect your voltage-measuring cables to the “V- Ω ” and the “COM” jacks.

3.3 Equipotential and Field Lines

- Connect two cables to the output of the power supply (black and red poles) and two of the yellow-metallic pins as described in the last step of the previous part. These are your output pins.
- Take one additional cable and connect one end to the black pole of the multimeter and the other to a different metallic pin. This is your reference or fixed-point probe. Take the pen-like pin and plug it into the red pole of the multimeter. This is your measuring probe.

- Take an empty sheet and put the two output pins on it to act as point-like sources. Push the pins in only a little bit and not all the way through.
- Insert the reference probe at an arbitrary position and move the measuring probe until the multimeter shows a value of zero. This means that the two points are at the same potential. Mark this position with one of the red markers provided, so as not to confuse it with the probe. Search for more equipotential points until you have enough to draw an equipotential contour.
- Move the reference pin to a new position and construct another equipotential line.
- Repeat for a total of about 5 equipotential lines.
- Given these equipotential lines, draw about 5 field lines using the yellow pen. Remember that the field lines and equipotential lines always intersect perpendicularly.
- Do your equipotential and field lines show the symmetries you would expect for this system? Just what symmetries does/should this system have?
- Obtain the equipotential and field lines for the parallel-plate configuration as well. (Stick the two pins from the function generator into the metal paint-covered part of the paper, after the paint has dried!).
- Are the equipotential lines as you expect them?
- What the major problems and uncertainties in this part?

3.4 Electric Shielding Theorem

Take the sheet with the conducting loop on it, and set up an electric field. You need the reference pin as in the previous part. Put it somewhere either within the loop or on the rim of it.

Put the reference pin anywhere within the closed surface or on its rim, and check that any point within the loop has the same potential as the reference pin. This demonstrates that all points within the closed loop have the same potential.

- What do you conclude about the shielding theorem? Does it hold or not?
- What would happen with your readings if you had drawn the loop badly such that the loop was not perfectly closed?

3.5 Linear Increase in Potential

The purpose is to verify that the potential increases linearly with distance as you move the probe from one end to the other.

Take the small strip and connect the two output pins to the metal covered ends. This produces a potential difference between the ends. Take the reference probe and remove the pin from the end. (You only need the cable and not the pin.) Connect it to one of the output pins. Take the measurement probe and measure how the voltage increases as you move further away from the reference point.

Plot a voltage vs. distance graph. Draw a best-fit line.

- Do the points fall on a reasonably straight line?
- How do you interpret the slope and intercept of this line?
- Is the field between the two poles homogeneous?

4 Applications

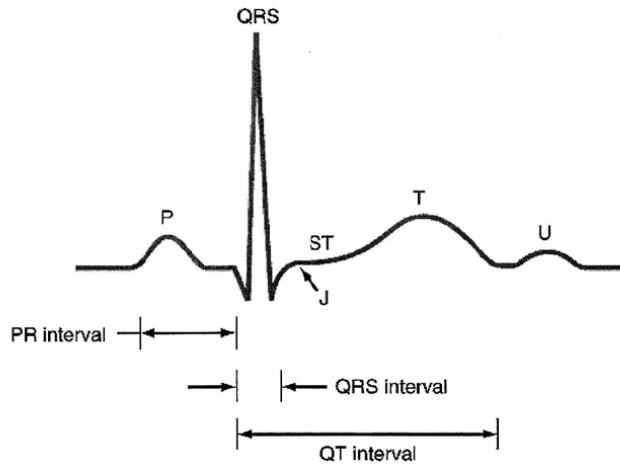


Figure 5: Sample Voltage Pattern for ECG

The canonical example in medicine for measuring voltage potentials and displaying them (with an oscilloscope) are ECG (Electrocardiography) and EEG (Electroencephalography). The basic idea of these two standard devices is fairly simple: by measuring the potential difference in time between two points (or usually several pairs of points) one gets a nice direct insight into the activity of the heart (brain) and can therefore detect dysfunctions easily. Let us, for example, take a closer look at how an ECG

works: If you record the electrical potential between two points along your chest you can record the voltage pattern shown in figure 5 (in time).¹

As you can see, the pattern observed can be split into different parts: At first the atrial cells are depolarized, giving the first signal (P wave). (This signal is relatively weak due to the small mass of the atrium.) After a delay one gets the QRS complex, which indicates the depolarization (wave) of the ventricles. After another delay one gets the T wave, which comes from the repolarization of the heart.² What do all these results, obtained from the heart as a total, mean in terms of processes going on in the single cells? The next figure 6 shows the measurement by your ECG in comparison to the potential in a single cell in a ventricle (obtained using a different method). First we recall that the interior and the exterior of the cell in their standard state have different ion concentrations and are therefore at different electrical potentials. (That is the zero line in the graph.) In the depolarization phase the ion channels in the cell membrane open and a flow of K^+ ions changes the potential inside the cell very rapidly. After this rapid change in potential one gets a plateau, which is mainly due to the inflow of Ca^{++} ions into the cell. (The Ca^{++} ions are much bigger than the K^+ ions, because they have a much larger hydrogen cover surrounding them and they diffuse much slower through the ion channels.) Finally in the repolarization phase the ion channels close again and the ion pumps in the cell membrane reestablish the initial ion concentrations.

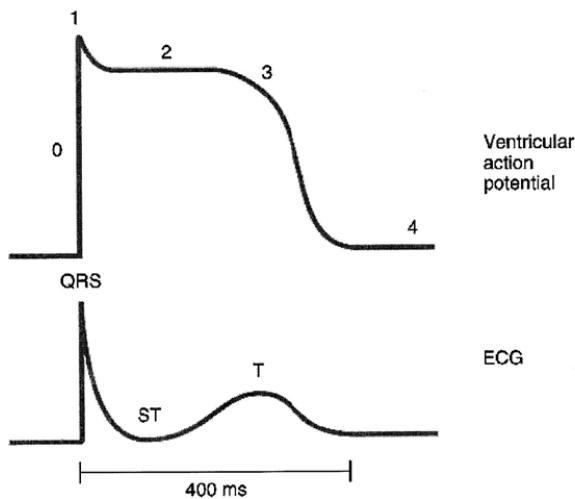


Figure 6: Comparison with Ventricular Potential

How can you now take advantage of this method for a diagnosis? First of all you

¹Remember: The ECG only shows the electrical polarization of the heart muscle. It does not show the contraction of the heart muscle!

²The interpretation of the U wave is still not yet 100% understood

don't want to take the measurements only along one single axis or plane, since e.g. if an infarction occurs on the front or back wall of your heart you are probably going to miss it. These days the usual way to get a 3-D picture of the position of the heart axis³ is obtained by measuring with multiple channels simultaneously between the points indicated in figure 7. (The contacts with an R are placed on the patients back.)

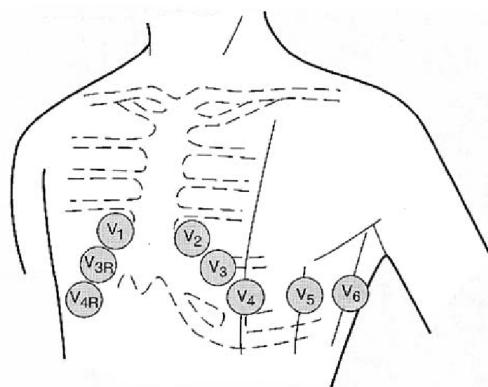


Figure 7: Channels of Measurement

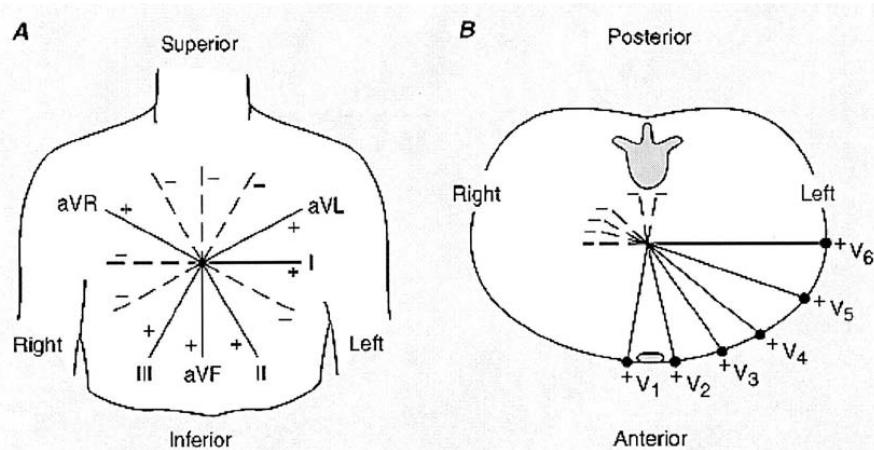


Figure 8: Electric Potentials and Information

From the location and amplitude of the main vector (axis) one can see for example if the muscle mass is increased on one side of the heart (hypertrophy). In that case the

³The heart axis is a simplified concept of the locations of the electrical potentials in the heart. One can think of the heart axis as a vector symbolizing the (physical) axis of the heart.

main vector is tilted.⁴ Another thing to look for is if the depolarization and repolarization was performed properly. For example if the ion channels in a certain region of the heart are destroyed by an infarction, then the electric potential between the depolarization and repolarization phase does not reach the zero level. By looking at your data you can not only locate the infarction, but also read off additional information, e.g. if the infarction cilled the tissue through all of the heart wall or only parts of it. (That determines your treatment of the patient!)

You can see that you can get a lot of useful information if you look at electrical potentials (and their change in time), as shown in figure 8.

Textbook references:

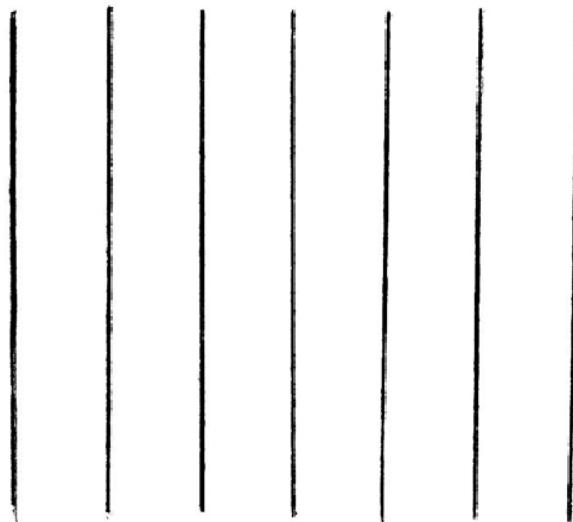
- Stein: *Internal Medicine*
- Harrison's *Principles of Internal Medicine*

⁴Also in pregnant women the heart as a total is slightly repositioned and therefore the electrical axis is a little bit off.

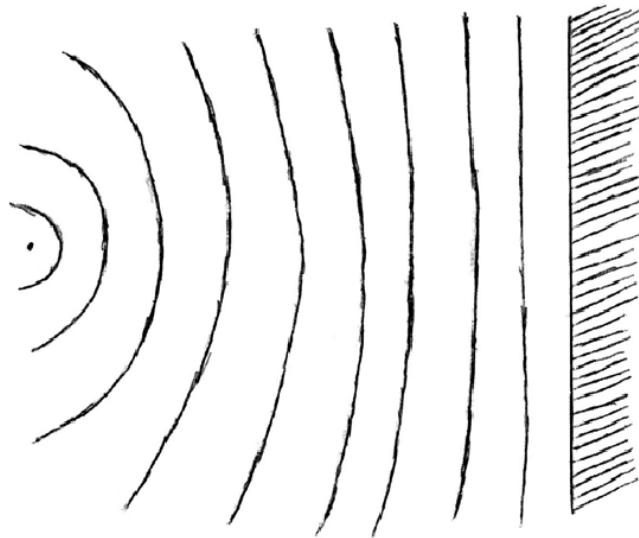
5 Lab Preparation Examples

Field Lines: Given the following equipotential lines, draw 5-10 field lines on each diagram.

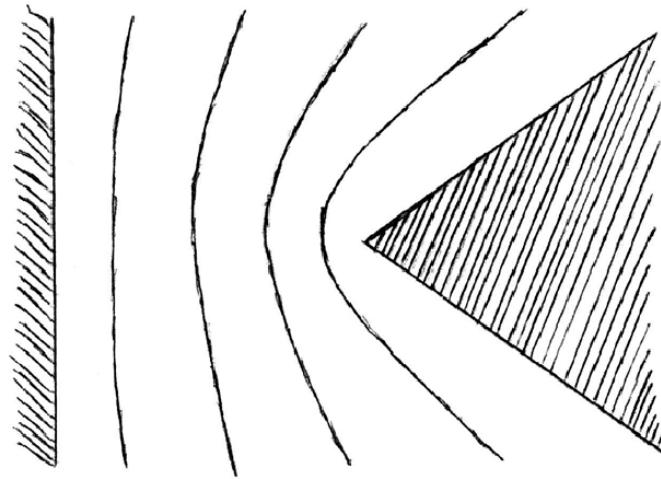
1.



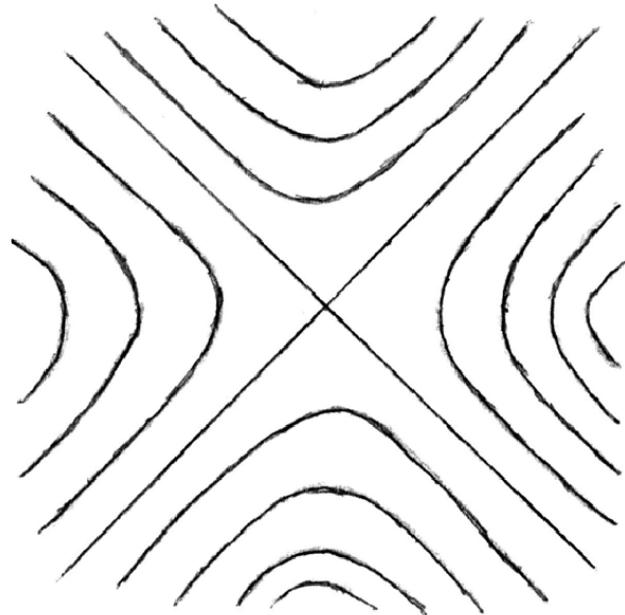
2.



3.



4.

Shield Theorem

5. Explain in a few sentences why you might be safe inside your new Volkswagen Beetle even when struck by a bolt of lightning.

Experiment 2-2

DC Circuits

Experiment 2-3

The Oscilloscope

1 Introduction

The cathode ray oscilloscope, one of the most useful tools in modern experimental physics, can measure potential differences which change rapidly with time – too rapidly to be followed by the needle of a simple voltmeter (some examples are shown in figure 1 below). You are encouraged to explore the operation of the oscilloscope by manipulating all of the controls.

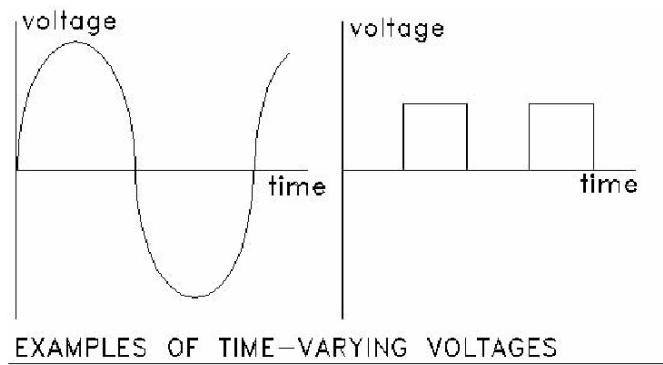


Figure 1: Examples of Time-varying Voltages

1.1 The Cathode Ray Tube

Figure 2 shows schematically the essential features of the cathode ray tube in the oscilloscope. The cathode ray tube is the fundamental component of the oscilloscope; in it the electron beam is deflected based on the potential difference that we are trying to measure. The electron beam is then incident on a fluorescent screen which allows us to see the results of the deflection caused by the potential difference.

Electrons leave the **heated cathode**, are accelerated through a fixed voltage, and emerge as a narrow beam focused through a hole in the **accelerating plate**. When the **electron beam** strikes the **fluorescent screen** on the face of the tube, it produces a small luminous spot. Between the accelerating plate and the screen are two pairs of parallel deflection plates – vertical and horizontal. A potential difference can be measured by applying it across one of these pairs of deflection plates. Before the potential difference is applied to either of these plate pairs it passes through a **vertical** or **horizontal amplifier** to increase the voltage of the signal. The applied potential

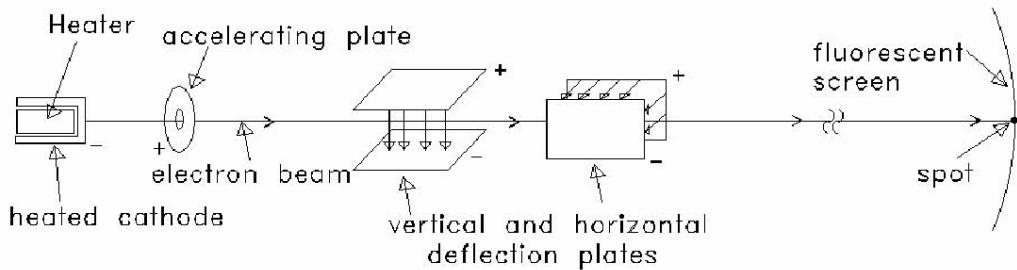


Figure 2: Schematic Diagram of a Cathode Ray Tube

difference results in a transverse uniform electric field between the plates, by which the beam is then deflected. The direction of the deflection (up-down or left-right) is determined by which of the two pairs of deflecting plates (vertical or horizontal) the voltage is applied to. Usually, the potential difference we are interested in measuring is applied to the vertical deflection plates.

The applied potential is directly proportional to:

1. the electric field strength between the plates ($V = E \cdot d$)
2. the transverse force on the electrons ($E = F/q$)
3. the resultant transverse acceleration of the electrons ($F = ma$)

and since the time to travel the length of the plates does not depend on the transverse force, it is also directly proportional to:

4. the net transverse displacement of the electrons $\left(x = \frac{1}{2}at^2; t = \text{const}, v_0 = 0 \right)$
5. the transverse component of their velocity ($v = at; t = \text{const}, v_0 = 0$)
6. the tangent of the angle at which the beam leaves the region between the plates $\left(\tan \theta = \frac{v_{\text{transverse}}}{v_{\text{parallel}}}; v_{\text{parallel}} = \text{const} \right)$

Therefore, when any of these quantities are multiplied by any factor, all of the others are multiplied by the same factor.

The net displacement of the spot on the screen will therefore be proportional to the applied potential difference. Therefore, by viewing the image on the oscilloscope, one can determine the voltage applied to the plates.

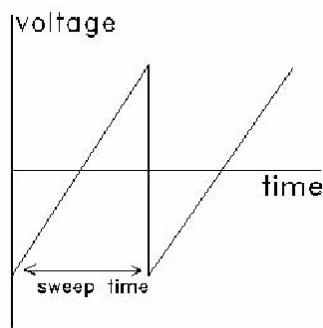


Figure 3: Horizontal Sweep Voltage

1.2 Display of Time-Varying Potential Difference

In order to display the time variation of a potential difference which is applied across the **vertical deflection plates**, the beam can be deflected horizontally by an internally-generated voltage which increases uniformly with time. Figure 3 indicates the “saw-tooth” voltage which, when applied to the **horizontal deflection plates**, sweeps the beam *horizontally* across the screen at *constant velocity*, and returns the beam to its initial position and repeats the sweep, etc.

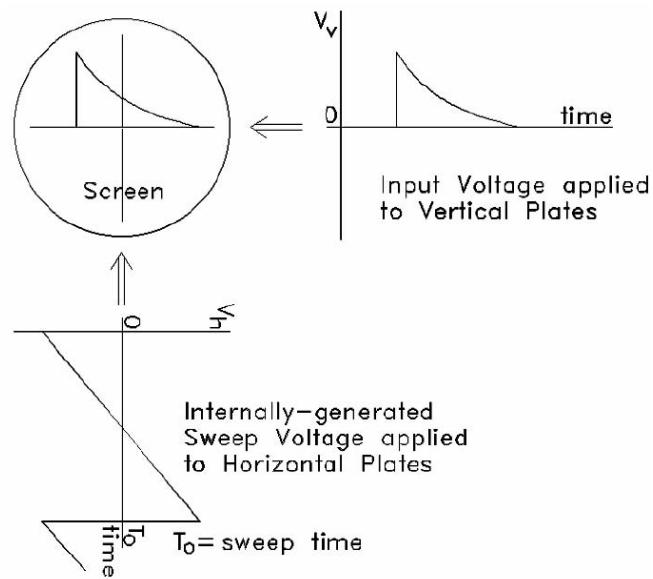


Figure 4: The image on the screen is formed by combining the constant horizontal sweep with a vertical deflection proportional to some time-dependent input voltage which you want to observe.

The process is somewhat similar to writing on paper. To write, one “wiggles” the pen up and down while moving the hand horizontally across the paper. Here the time-varying voltage applied to the vertical plates “wiggles” the electron beam up and down while the linear sweep voltage moves the beam horizontally with constant velocity (See figure 4). The beam leaves its written trace on the tube screen. The motion in the vertical plane therefore usually corresponds to the voltage we are measuring, while the motion in the horizontal plane usually corresponds to the passage of time.

Because of the persistence of vision (approximately 0.05 sec) and because of the screen fluorescence, we see a plot of the potential as a function of time displayed on the screen.

Usually one observes the repeating pattern of a periodically varying voltage by appropriate adjustment of the sweep time and synchronization of the start of successive sweeps. Oscilloscope controls allow one to determine when the sweep begins (on the oscilloscope this is called the “trigger”). See the description of the horizontal sweep trigger controls (17-21) in the instructions for the oscilloscope in the next section.

2 The Oscilloscope

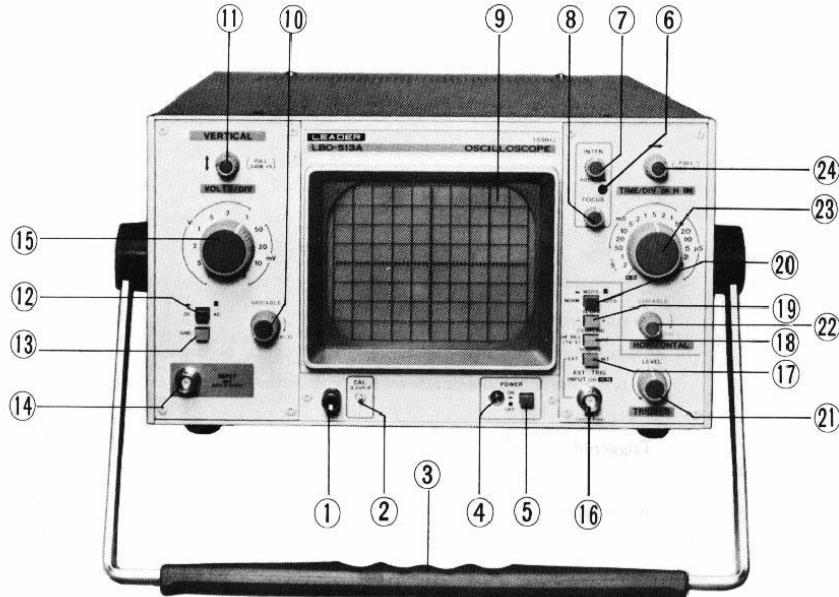


Figure 5: Leader Model LBO-513A Oscilloscope

The front panel of the Leader Model LBO - 513A oscilloscope is shown in the figure 5. The inputs and controls with which you will be concerned can be divided into three groups depending on the part of the oscilloscope circuitry with which they are associated. All oscilloscopes include this basic set of controls. The numbers in the diagram refer to the controls described below.

2.1 FORMATION OF THE ELECTRON BEAM

(5) POWER

Push in to turn the power on.

(7) BEAM INTENSITY

The **INTEN** control adjusts the intensity or brightness of the trace. Clockwise rotation increases the intensity.

NOTE: When the spot on the screen is stationary, keep the intensity very low in order not to damage the fluorescent screen at that point. In general, do not leave the intensity higher than necessary for reasonable visibility.

(8) BEAM FOCUS

The **FOCUS** control adjusts the sharpness of the electron spot on the screen.

2.2 VERTICAL DEFLECTION OF THE ELECTRON BEAM

(10) VERTICAL AXIS SENSITIVITY: VARIABLE FINE ADJUSTMENT.

The vertical axis **VARIABLE** control is used for fine adjustments of the vertical axis sensitivity. The sensitivity is reduced by turning the vertical axis **VARIABLE** knob counterclockwise.

In the fully clockwise position (with the knob clicked into this position) quantitative data about the potential difference can be read from the **VOLTS/DIV** setting.

(15) VERTICAL AXIS SENSITIVITY: RANGE SELECTION IN VOLTS/DIV.

The **VOLTS/DIV** 11-position switch determines the sensitivity of the vertical amplifier. The 11 ranges are indicated in volts per division on the front panel.

The indicated sensitivities are only correct if the vertical axis **VARIABLE** control is in the calibrated position (fully clockwise – see (10)) and the **VERTICAL** position control is pushed in (the Gain = 1 position – see (11-B)).

(11-A) VERTICAL POSITION:

This control is used to move the displayed trace up or down. Turning the knob clockwise moves the trace upward.

(11-B) VERTICAL AXIS MAGNIFIER: GAIN × 5:

If the **VERTICAL** position control is pulled out, the vertical axis sensitivity is increased by a factor of 5.

(12) AC-DC INPUT SELECTOR:

AC stands for “alternating current,” a current that changes with time, while DC stands for “direct current,” a current that is constant in time.

With the button in (DC position), the input terminal is directly coupled to the vertical amplifier and both AC and DC signals are shown. With the button out (AC position) the DC signal is blocked by a capacitor so only the AC, or varying, signal is shown while the DC, or constant, signal is subtracted off.

(13) GND GROUND SWITCH:

With the **GND** button in, the input to the vertical amplifier is grounded and

therefore there is no deflection in the vertical direction.

(14) **VERTICAL INPUT:**

This is the input terminal for the vertical amplifier. The maximum permissible input voltage is 600 Volts.

2.3 HORIZONTAL DEFLECTION OF THE ELECTRON BEAM

(16) **EXT TRIG INPUT or H IN:** EXTERNAL TRIGGER INPUT OR HORIZONTAL INPUT:

This input is used for externally triggering the horizontal sweep or for the HORIZONTAL INPUT in an x-y plot. This input is DC coupled (either a constant or a time varying voltage may be used – compare with (12) above) and the maximum allowable applied voltage is 100 volts.

(17-21) **HORIZONTAL SWEEP TRIGGER CONTROLS.** There are four switches and one knob that are used in determining the conditions for triggering the horizontal sweep which starts the electron beam traveling horizontally across the screen

(17) **SOURCE = INT** Sweep is triggered on the internal signal in conjunction with the **VERTICAL INPUT**.

SOURCE = EXT Sweep is triggered on an external signal fed into the **H IN HORIZONTAL INPUT** (16).

(18) **COUPLING = AC** Sweep is triggered on any AC signal.

COUPLING = HF-REJ(TV-V) Sweep is triggered only on signals with a frequency below 20 kHz.

(19) **SLOPE = +**Sweep is triggered on the positive slope of the signal.

SLOPE = -Sweep is triggered on the negative slope of the signal.

(20) **MODE = AUTO** Sweep is automatically triggered regardless of amplitude or frequency.

MODE = NORM Sweep is triggered when the signal is greater than the trigger level set by the **LEVEL** control.

(21) **TRIGGERING LEVEL** This control determines the voltage level that will trigger the horizontal sweep when the **MODE = NORM** is selected.

(22) **HORIZONTAL SWEEP TIME: VARIABLE FINE ADJUSTMENT:**

The horizontal axis **VARIABLE** control is used for fine adjustments of the horizontal sweep time per division. The sweep time per division is increased by turning the horizontal axis **VARIABLE** knob counterclockwise. Only in the fully clockwise

position (with the knob clicked into position) can the sweep time be read from the **TIME/DIV** setting.

(23) HORIZONTAL SWEEP TIME: RANGE SELECTION IN **TIME/DIV** or **H IN** HORIZONTAL INPUT SELECTOR:

This control is used to select the horizontal sweep time per division on the scope face. The indicated sweep times per division are only correct if the horizontal axis **VARIABLE** control is in the calibrated position (fully clockwise – see (22)) and the **HORIZONTAL** position control is pushed in (no magnification position – see (24-B)). In the fully counterclockwise position the sweep oscillator is disconnected and the signal on the **H IN** terminal is applied to the horizontal plates.

(24-A) HORIZONTAL POSITION:

This control is used to move the trace left or right. Turning the knob clockwise moves the trace to the right.

(24-B) HORIZONTAL SWEEP TIME MAGNIFIER: **MAG × 5**:

If the **HORIZONTAL** position control is pulled out, the horizontal sweep time per division is reduced by a factor of 5 and the wave form displayed is consequently magnified by a factor of 5.

3 Procedure

3.1 Check the linearity between beam deflection and applied voltage

Use the slide wire resistor, connected in series with a fixed DC power supply, to serve as a variable voltage source. Make sure that the slide wire resistance is connected in a closed circuit to the DC power supply and then connect (14) VERTICAL INPUT with two cables to the slide wire resistance. One cable should be connected to the fixed point on one end of the slide wire resistor and the other should be connected to the movable slider.

Obtain a spot – *at low intensity* – on the screen by disconnecting the horizontal sweep voltage from the horizontal plates. The “sweep voltage” is applied to the horizontal plates only when (23) is set on one of the “sweep time” settings, so in this case, set (23) to the **H IN** position since we do not want the horizontal sweep voltage applied in this case.). Also set control (12) to DC, since we are measuring a potential difference that is constant in time. With your voltage divider, check to see if the beam deflection is proportional to the voltage applied to the oscilloscope input terminals. Present your results graphically.

3.2 Use the oscilloscope to display periodic and aperiodic signals

a) Use the signal generator to put sine waves and square waves of various frequencies on the vertical plates. Vary the amplitude of the input voltage, and vary the sweep times and the method of synchronizing the sweep.

b) Observe some irregularly varying potential such as the output of a microphone into which you speak, hum, whistle, etc. Determine the range of frequency of your voice by humming as low and as high a note as possible.

c) Use the signal generator to put a variable frequency sine wave on the vertical plates, and have a small transformer connected to the horizontal plates to furnish them with the 60 cycle/sec AC line signal. Vary the frequency and see whether you can observe a “Lissajous Figure” – the closed curve produced when the ratio of frequencies is the ratio of small integers.

3.3 Use the signal generator and the oscilloscope to measure the gain of an audio amplifier as a function of frequency

A high fidelity amplifier is one that has constant amplification over the normal range of audio frequencies. The amplification (or “gain”) is defined as the ratio of the peak-to-peak voltage of the output signal (delivered to a loud speaker) to that of the input signal. Normally, of course, the input signal comes from a microphone, a record player, or a radio tuner. Here, you can use the signal generator to produce sine waves of variable frequency and amplitude.

It is convenient to keep the amplitude of the input signal constant; it is important to keep the amplitude small so that the amplifier does not “saturate” and distort the shape of the output wave.

Plot a curve of gain versus frequency for a few settings of the “tone” control. A curve such as this is the standard criterion for evaluating (and advertising) a high-fidelity amplifier.

If you have a (relatively small) amplifier which you wish to test, you may bring it to the laboratory to use instead of the small (and relatively low fidelity) amplifiers provided.

Experiment 2-4

Capacitance

1 Introduction

A capacitor is a device for storing electric charge and energy. For simplicity, an ideal capacitor can be considered as a pair of parallel conducting metal plates, as shown in figure 1. When a charge $+Q$ is placed on the upper plate and $-Q$ on the lower plate, a potential difference V is established between the plates, and the quantities Q and V are related by the expression:

$$Q = CV \quad (1)$$

where the capacitance C is determined by the size and separation of the plates.

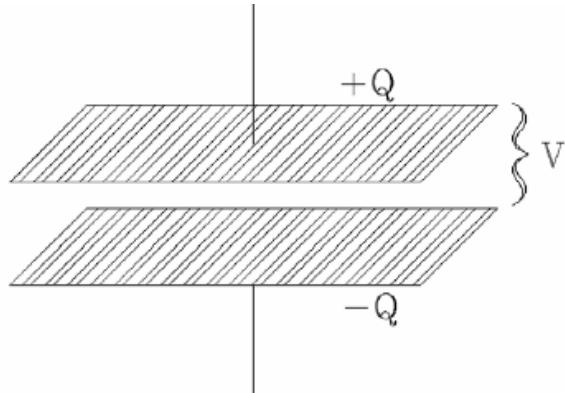


Figure 1: An Ideal Capacitor

1.1 Discharging a Capacitor

If a wire is attached to the upper plate and then touched to the lower plate, charge will flow through the wire until the charge q and potential difference V are zero. (We will use Q to indicate the initial charge and q for the time dependent charge.) This process, known as discharging the capacitor, does not occur instantaneously. Instead, as the charge flows from one plate to the other, the potential difference decreases, and the current (or rate of change of charge) gradually diminishes. Since the rate of change is proportional to the amount of charge remaining, the current is initially large and then *decays exponentially* with increasing time. Exponential growth or decay occurs whenever the rate of change of a quantity is proportional to the quantity present at that time. Exponentials are also characteristic of the growth rate of cells or organisms and the decay of radioactive isotopes.

The changing current which discharges a capacitor can be calculated as a function of time by considering the circuit shown in figure 2. The capacitor C is initially charged by a battery of emf ε , placing a charge $Q = C\varepsilon$ on the plates. When the switch is closed, a current I flows in the circuit, and according to Ohm's Law the voltage drop across the resistance is IR .

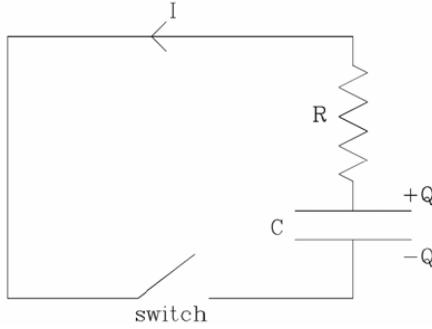


Figure 2: Discharging Circuit of a Capacitor

The resistance may be a separate circuit element or merely the resistance of the wire itself. Since the sum of the voltage drops around the circuit must be zero, we obtain the equation:

$$IR + \frac{Q}{C} = 0 \quad (2)$$

The equation can be rewritten in terms of charge alone from the definition of current, $I = dQ/dt$ ($= \lim_{\Delta t \rightarrow 0} \frac{\Delta Q}{\Delta t}$),

$$R \frac{dQ}{dt} + \frac{Q}{C} = 0 \quad (3)$$

Equation (3) is a differential equation for Q and has a solution which gives the time dependence of the charge

$$Q(t) = C\varepsilon e^{-t/RC} \quad (4)$$

Students familiar with calculus can verify this result by differentiating equation (4) and combining $Q(t)$ and dQ/dt to show that equation (3) is satisfied. From the definition of current, $I = dQ/dt$, we find that

$$I(t) = -\frac{\varepsilon}{R} e^{-t/RC} \quad (5)$$

Note that at $t = 0$, when the switch is just closed, $I = -\varepsilon/R$. As time increases, the current $I(t)$ gets smaller and reaches zero as time goes to infinity.

1.2 Charging a Capacitor

The process of charging an uncharged capacitor has many similarities with the process of discharging described above. In this case, a battery with an emf of ε volts is connected in series with a resistance R and the capacitance C , as indicated in figure 3.

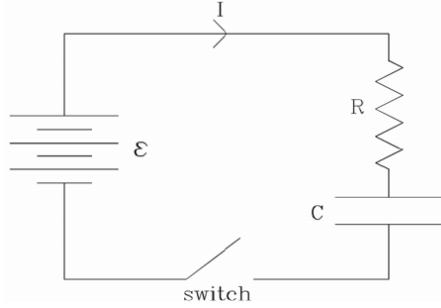


Figure 3: Charging Circuit of a Capacitor

When the switch is first closed, there is no charge on the capacitor and thus no voltage across it, and the full voltage ε appears across R , causing maximum current I to flow. As the capacitor becomes charged, the voltage-drop IR (and thus the current I) gradually decreases. I approaches zero as the capacitor is charged toward its maximum potential difference of $\varepsilon = Q/C$.

The equation which indicates the sum of voltage drops around the circuit loop at any time during the charging is:

$$\varepsilon = IR + \frac{Q}{C} \quad (6)$$

Again, this can be converted to a *differential equation*, the solution of which is:

$$I(t) = \frac{\varepsilon}{R} e^{-t/RC} \quad (7)$$

Comparing equation (5) with (7), we note that the current for discharging a capacitor from a given potential ε decreases in time identically to the current for charging the capacitor through the same resistance to the same final voltage ε . The difference in the sign of I indicates, of course, that the two currents flow in opposite directions.

1.3 Graphical Presentation of Exponential Decay

In equations (5) and (7), the symbol e stands for the constant $2.71828\cdots$, the *base of natural logarithms*. A graph of the function $y = ae^{-x/b}$ is given in figure 4, where $y = a$ at $x = 0$ and then decays exponentially as x increases, reaching a value of $y = ae^{-1} = 0.37a$ at $x = b$.

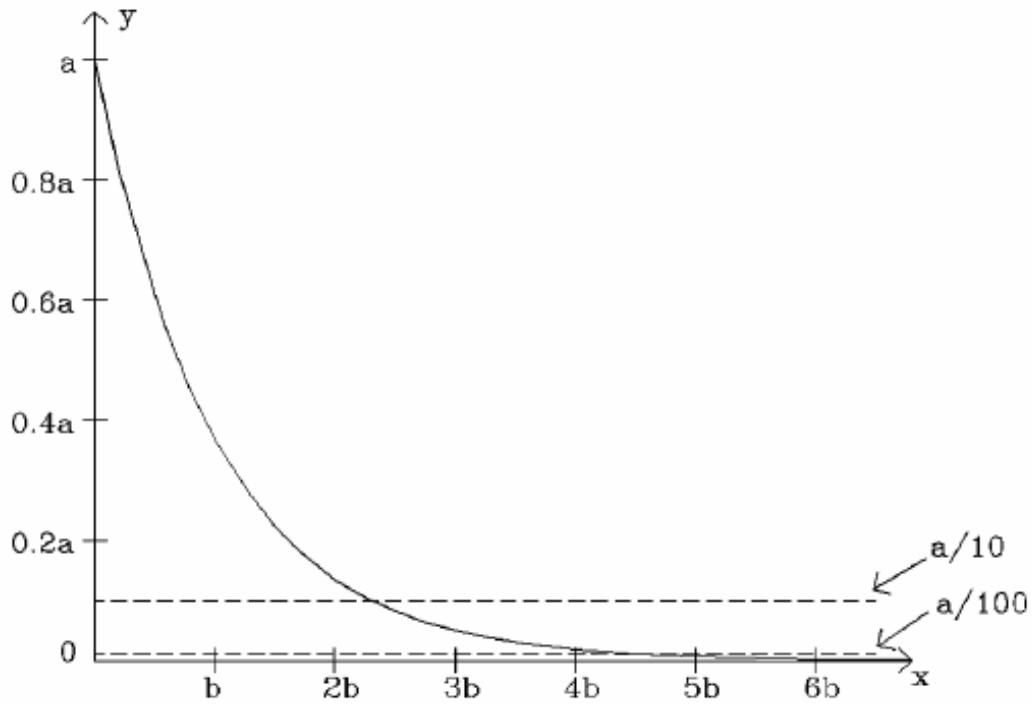


Figure 4: Graph of $y = ae^{-x/b}$

Although it is possible to determine the values of a and b from experimental results plotted on such a curve, it is more informative to plot $\log y$ versus x , or to use *semi-log* graph paper, where the horizontal lines and scales of the y -axis are spaced in equal increments of $\log y$. Then, taking the logarithm (base 10) of $y = ae^{-x/b}$ yields:

$$\log y = \log a - \frac{x}{b} \log e = \log a - \left(\frac{\log e}{b} \right) x = \log a - \left(\frac{0.434}{b} \right) x \quad (8)$$

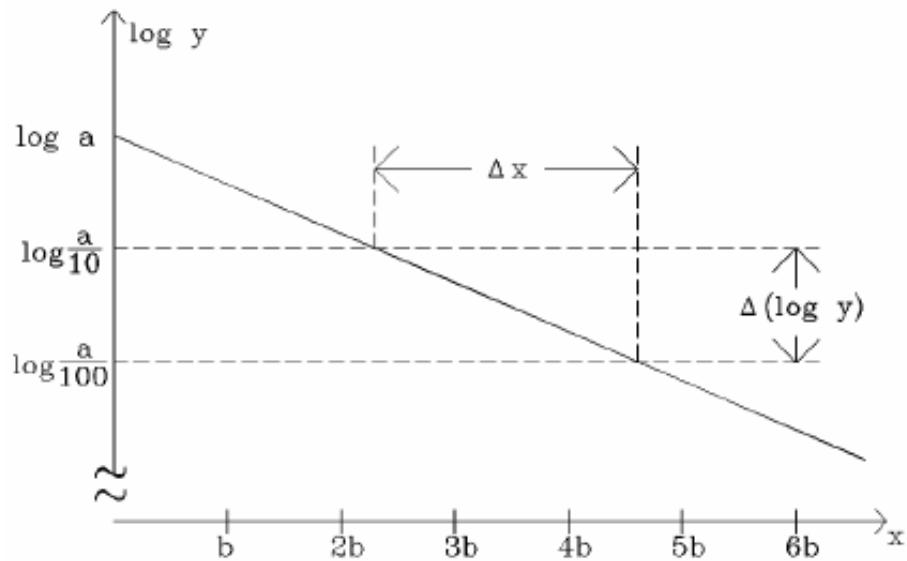
which is a linear equation in x representing a straight line with slope $-0.434/b$, if $\log y$ is plotted against x as in figure 5.

A simple way to determine b from a straight line drawn through the experimental points on such a plot is to note the values of x at which the line crosses values of y differing by a factor of 10. Then, using the values shown in figure 5,

$$\Delta(\log y) = \log \frac{a}{100} - \log \frac{a}{10} = (\log a - \log 100) - (\log a - \log 10) = -2 + 1 = -1 \quad (9)$$

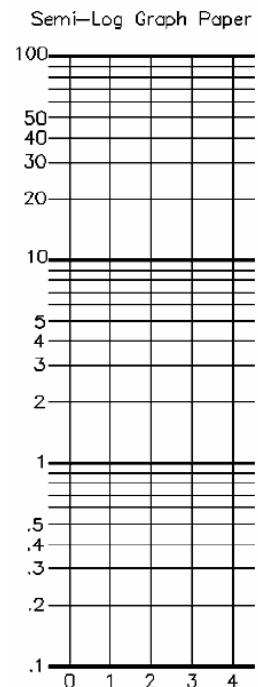
Since from equation (8) the slope $\Delta(\log y)/\Delta x$ equals $-0.434/b$, therefore

$$\frac{-1}{\Delta x} = \frac{-0.434}{b}, \quad \text{or} \quad b = 0.434 \Delta x \quad (10)$$

Figure 5: Plot of $\log y$ versus x

In the present experiment, y represents the current I ; x is the time t ; a is the initial value of I ; and b is the time constant RC .

Note: You may produce a graph like the one in figure 5 by either: 1) Plotting $\log y$ on the y axis of normal graph paper; or 2) Plotting y on semi-log graph paper (shown in the right). Semi-log graph paper (provided in the lab) has the y axis already logarithmically compressed, so you do not need to calculate $\log y$ for each point – you can simply plot the y values themselves. With semi-log graphs you can verify quickly whether the data behaves exponentially or not; $y = 1/x$, for example, will not produce a straight line on a semi-log plot. Note also that on semi-log graphs y approaches 0 only exponentially – there is no $y = 0$.



2 Procedure

2.1 Large RC-charging

Figure 6 shows the circuit to be wired for measuring the current I as a function of time for charging a capacitor C through a resistance R to a voltage ε . For this part of the lab you will use the bank of three large capacitors which are attached to a piece of wood with three switches and the label “30 MFd”. (In this case the prefix “M” in MFd indicates micro, μ , or 10^{-6} .) Be sure to use a low-voltage (15 volts max.) power supply and to connect the positive terminal on the power supply to the positive terminal on the microammeter. Before closing the switch to the power supply and starting each new measurement, momentarily connect a low resistance across the capacitor in order to start with no charge on the capacitor plates.

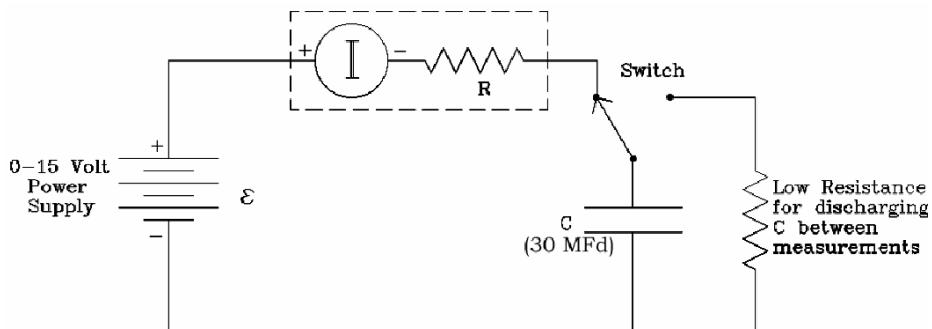


Figure 6: Large RC-charging Circuit

Measure I at a series of time intervals after closing the switch. Plot the results on semi-log paper and determine the RC time constant.

Repeat the above for a different value of C . (Note that the “capacitor” supplied is in fact a bank of capacitors with the provision for varying C by switching in different numbers of the capacitors in parallel).

2.2 Large RC-discharging

Use the circuit shown in figure 7 to measure the discharge current through the resistance R of the same capacitors C used in Part 1. You will use the same equipment as you used in Part 1 of the lab.

Before closing the switch to begin the measurement, momentarily connect the power supply directly across the capacitor in order to give it an initial charge. Then disconnect

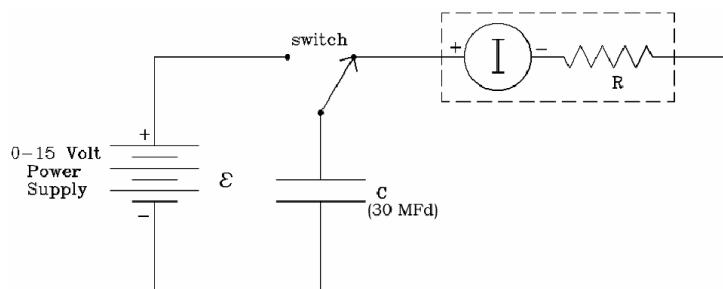


Figure 7: Large RC-discharging Circuit

the power supply, close the switch, and measure I as a function of time. Again, plot the results on semi-log paper and determine RC . Compare with the results found in Part 1 for the same values of C .

2.3 A Relaxation Oscillator

A neon bulb has the property of having a very high resistance (almost infinite) until the voltage applied to it is high enough to “break down” the gas, at which point the bulb lights and its resistance becomes very low. In this part of the lab you will use a high voltage power supply. Thus, *it is essential that you are sure to use the small capacitor with a value of $0.082 \mu\text{F}$* rather than the bank of capacitors you use in the first two parts of the lab. Wire a neon bulb in parallel with the capacitor in the circuit of figure 8.

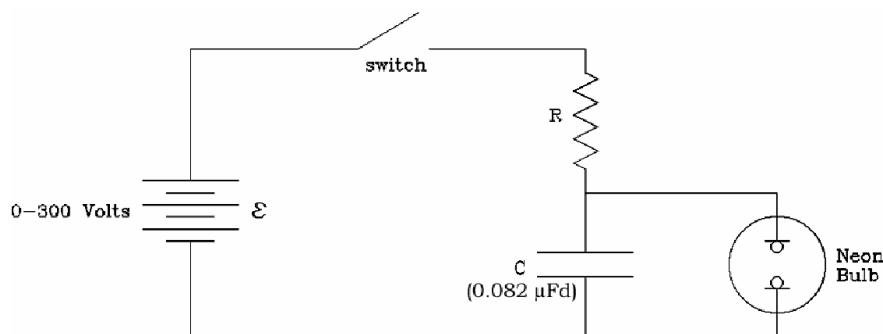


Figure 8: Relaxation Oscillator Circuit

CAUTION: Since a relatively high-voltage power supply is necessary for this part of the experiment, it is important not to touch any exposed metal parts of the circuit while the power supply is connected to the circuit and turned on, whether or not the

switch is closed. Since the capacitor stores charge, *it may be charged even if the voltage supply has been removed*. Be sure to discharge the capacitor fully by simultaneously touching an insulated wire to each end of the capacitor before touching any of the metal parts of the circuit.

When the switch is closed, the capacitor will start to be charged as in Part 1, with the time constant RC , and the high resistance of the neon bulb will have negligible effect on the circuit. However, when the capacitor is charged to sufficiently high voltage, the neon bulb will light. The capacitor will then be quickly discharged through the bulb. If R is sufficiently large, there will not be sufficient current to keep the bulb lit after the capacitor is discharged. The bulb will then be extinguished; it will return to a state of high resistance, and the charging process will start again. For a given applied voltage ε , and a neon bulb with a given breakdown voltage, the period of this repetitive “oscillation” will thus be determined by the value of RC .

Measure the period for different combinations of R and C , and verify its dependence on the product RC .

2.4 Demonstration of Short RC Time (to be set up by Lab Instructor)

If the RC time constant of a circuit is too short to be detected using an ammeter, the effect can be displayed on an oscilloscope as illustrated in figure 9. In this case, it is convenient to use a square-wave generator, which switches a voltage on and off at a variable repetition rate.

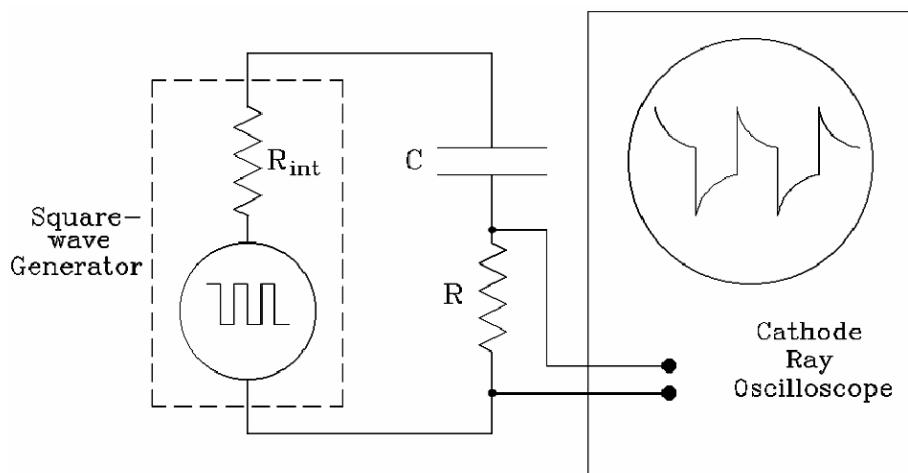


Figure 9: Using an Oscilloscope to Study the Time-dependence of Voltage

Use the smallest values of R and C available. From observation of the oscilloscope trace, plot I versus time on semi-log graph paper and determine the value of RC . Compare the result with the value of RC calculated from the labels on the circuit elements. (Remember that the total R of the circuit includes any internal resistance in the square-wave generator).

Experiment 2-5

The Magnetic Field

OBJECT: To determine the strength of the magnetic field in the gap of an electromagnet by measuring: (I) the force on a current-carrying rod, (II) the effect of the change of magnetic flux through a coil when the coil is inserted into or removed from the region containing the magnetic field.

CAUTION: (1) Always reduce the current through the electromagnet to zero before opening the circuit of the magnet coils. (2) Remove wrist watches before placing hands near the magnet gaps.

1 Part I - Force on a Current-carrying Wire

1.1 Physical Principles

Consider a rigid length of wire rod L , held horizontally and normal to the direction of a uniform, horizontal magnetic field B . If a current i is passed through the wire, as indicated in Figure 1, then there will be a vertical force $F = iLB$ on the wire rod.

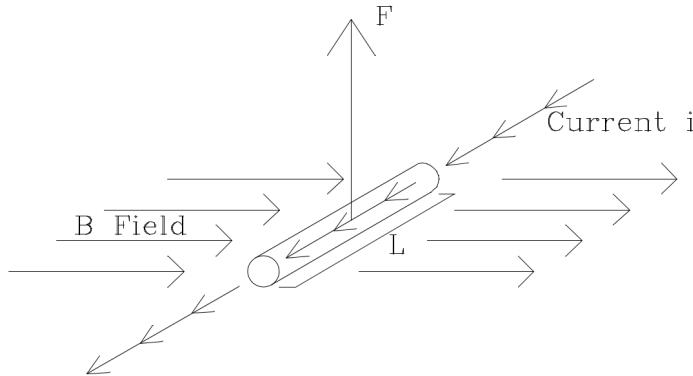


Figure 1: Force on a Current-carrying Wire

1.2 Experimental Apparatus

The experimental set-up is shown in Figure 2. The horizontal magnetic field B is produced in the air gap of a “C-shaped” iron electromagnet. The strength of B is determined by I , the current in the magnet coils, which is supplied by an adjustable low-voltage power supply.

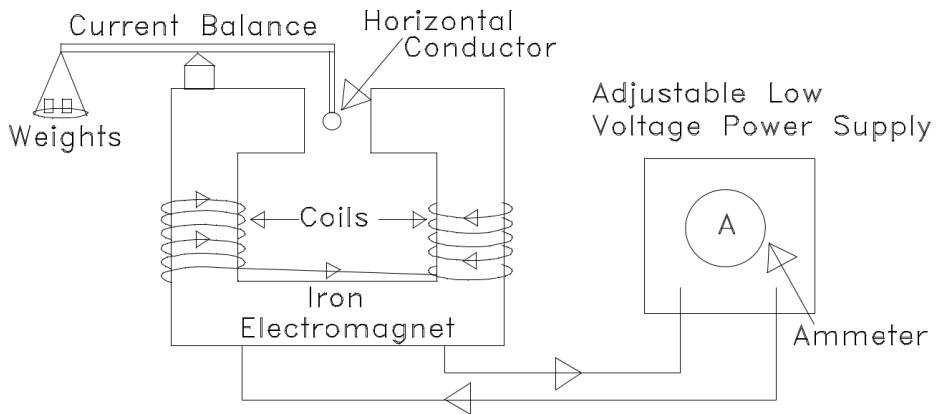


Figure 2: Experimental Set-up

A more detailed drawing of the balance and electro-magnet arrangement is shown in Figure 3.

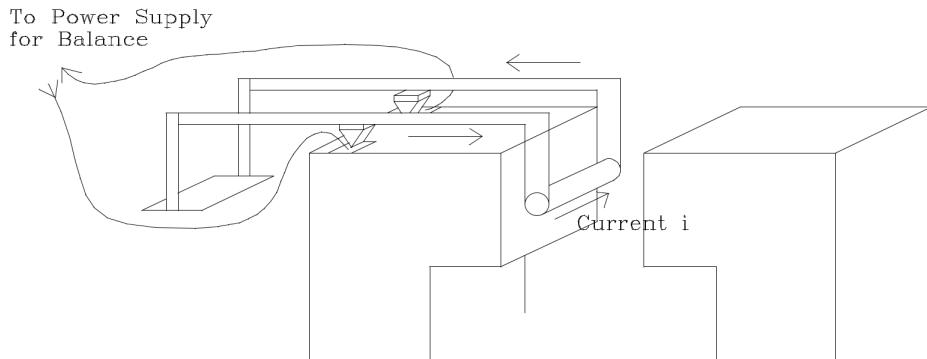


Figure 3: Setup of the Wire and Balance for Force Measurement

Note that the current i through the horizontal conductor is supplied by a separate power supply. The current balance is constructed out of conducting and insulating materials such that current can enter through one side of the knife edge, flow through one side of the balance arm to the horizontal conductor, and then flow back through the other side of the balance arm and out through the other knife edge.

The current i in the balance is provided by the HP E3610A power supply, which can operate either in constant voltage or constant current mode. The voltage dial sets the *maximum voltage* the device will supply to the circuit; the current dial likewise sets the *maximum current*, and if the circuit tries to draw more current, the power supply will reduce the voltage until it reaches whatever value is needed to maintain the maximum current (by $V = IR$).

To use the power supply in constant current mode, begin with the current dial

turned all the way down (counter-clockwise) and the voltage dial turned all the way up (clockwise). Set the range to 3 Amps, and connect the leads to the + and – terminals. You can now set the current to the desired level—note that the digital meters on the power supply show the *actual* voltage and current being supplied, so you will not normally see any current unless the leads are connected to a complete circuit. (If you want to set the current level without closing the circuit, you can hold in the CC Set button while you turn the current dial.) Once you close the circuit, the voltage adjusts automatically to maintain the constant current level, and the CC (Constant Current) indicator light should be on.

1.3 Procedure

Set the current I through the electromagnet at 5 amperes. Then determine the weights needed to bring the balance to equilibrium for several different values of the current i through the balance. (Note that it is easier to make the final adjustment on i once weights have been selected for the approximate current value.) Present the measurements of force versus balance current graphically, and from the graph find a value for B . Repeat the above for magnet currents of 4 amperes, 3 amperes, and 2 amperes. Determine B in each case. All your results may be plotted on the same graph. Draw a diagram similar to Figure 3, and indicate the direction of i , F , and B for your set-up.

2 Part II - EMF Induced in a Moving Coil

2.1 Physical Principles

A more accurate method of measuring B is to use a **charge integrating circuit** to determine the total charge which flows through the circuit of a small search coil when it is suddenly inserted into (or removed from) the region of magnetic field. The EMF ε induced across the coil by a change of magnetic flux $\Delta\Phi$ through the coil is given by Faraday's Law:

$$\varepsilon = N \frac{\Delta\Phi}{\Delta t} \quad (1)$$

where N is the number of turns in the search coil.

If the total resistance of the coil and its circuit is R , then the resulting instantaneous current $i = \Delta Q / \Delta t$ flowing through the circuit will be

$$\frac{\Delta Q}{\Delta t} = i = \frac{\varepsilon}{R} = \frac{N}{R} \frac{\Delta\Phi}{\Delta t}. \quad (2)$$

The instantaneous induced EMF and current will thus depend on the time Δt , i.e. the speed with which the search coil is moved into the field. However, the total charge

which flows during the motion is independent of Δt :

$$\Delta Q = \frac{N}{R} \Delta \Phi. \quad (3)$$

If the search coil is started in a field-free region and is plunged into the magnet gap so that B is normal to the plane of the coil, then $\Delta \Phi = A_{\text{coil}}B$, where A_{coil} is the area of the coil, and therefore:

$$\Delta Q = \frac{N}{R} A_{\text{coil}} B. \quad (4)$$

Note that both the inner and outer diameter of the search coil should be measured, since the innermost windings have a different area than the outermost ones.

2.2 Experimental Apparatus

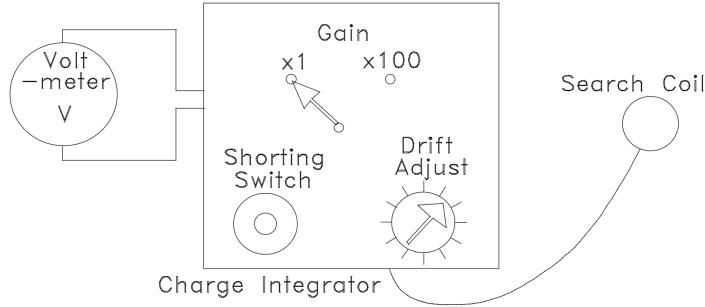


Figure 4: The Magnetic Field Module

A charge integrator (the Magnetic Field Module shown in Figure 4) is used to measure the ΔQ produced by the EMF induced in the search coil. A capacitor in the module stores the charge ΔQ , and the voltage across this capacitor (read on the external voltmeter shown) is proportional to ΔQ . Therefore

$$V = K \Delta Q = K \frac{N}{R} \Delta \Phi \quad (5)$$

where K is a constant that depends on the capacitance and gain of the integrator circuit. Instead of trying to calculate a value of KN/R in terms of the components, it is more direct to calibrate the combination of the search coil and the integrator circuit by measuring V for a known $\Delta \Phi$. This known magnetic flux can be created by passing a measured current I_{sol} through a long air-core solenoid of n turns per meter and with a cross-sectional area of A_{sol} so that

$$\Delta \Phi_{\text{sol}} = B_{\text{sol}} A_{\text{sol}} = \mu_0 n I_{\text{sol}} A_{\text{sol}} \quad (6)$$

where μ_0 is the permeability of free space ($\mu_0 = 4\pi \times 10^{-7} \text{ T} \cdot \text{m/A}$).

2.3 Procedure

Connect the apparatus as shown in Figure 4. Turn on the power supply. Before any measurements are made, depress the shorting switch, then release the shorting switch and turn the drift adjust control to minimize the drift in the output voltage as observed on your meter. The shorting switch must be used to discharge the integrating capacitor prior to each measurement. Also, the drift adjust setting should be checked occasionally. If the gain setting on the Magnetic Field Module is changed, the drift adjust control must be reset.

Magnetic field measurements are made by inserting or removing the search coil from the region containing the field to a field-free region or by leaving the search coil stationary and turning the field on or off.

Since the magnetic field in the large iron-core electromagnet is much greater than the field in an air-core solenoid, the Magnetic Field Module was designed with two gain settings. In the gain=100 position, the module is 100 times as sensitive as in the gain=1 position. For measuring fields generated by the air-core solenoid, set the gain at 100. For fields generated by the large electromagnet, set the gain at 1.

Connect the solenoid to the + and – terminals of the HP Power Supply. A three position (on-off-on) reversing switch is part of the solenoid circuit. Flip the switch to one of the on positions and adjust the power supply current such that two amperes is flowing through the solenoid.

Slide the search coil over the solenoid and, while holding the search coil at the center of the solenoid, discharge the Magnetic Field Module and then turn the current through the solenoid off using the three position switch (or turn the current from off to on). Take several readings and record the voltage on the integrator and the current through the solenoid. Then the result is

$$V_{\text{sol}} = 100 \cdot K \cdot \frac{N}{R} \cdot \Delta\Phi_{\text{sol}} = 100 \cdot K \cdot \frac{N}{R} (\mu_0 n I_{\text{sol}}) A_{\text{sol}} \quad (7)$$

Now set the gain to 1 on the Magnetic Field Module, and readjust the drift controls. Set the current for the large electromagnet to 5 amps, and use the search coil to measure the resulting B . Move the search coil *gently and smoothly* into the region of the magnetic field (do not move the coil hastily as you may damage it by striking against the magnet itself). Record V_{mag} . The corresponding equation is:

$$V_{\text{mag}} = K \cdot \frac{N}{R} A_{\text{coil}} B_{\text{mag}} \quad (8)$$

Combine the results of Eq(7) and Eq(8) to determine the value of B for the large electromagnet when the current is 5 amps, and then do the same for the other magnet currents of 4, 3, and 2 amps. Compare these values for B with those obtained in Part I.

Experiment 2-6

e/m of The Electron

1 General Discussion

The “discovery” of the electron by J. J. Thomson in 1897 refers to the experiment in which it was shown that “cathode rays” behave as beams of particles, all of which have the same ratio of charge to mass, e/m . Since that time, a number of methods have been devised for using electric and magnetic fields to make a precise measurement of e/m for the electron. When combined with the value of the electron’s charge, which is measured in the Millikan Oil Drop Experiment, the determination of e/m leads to an accurate value of the mass of the electron. In the present experiment, electrons are emitted at a very low velocity from a heated filament, then accelerated through an electrical potential V to a final velocity v , and finally bent in a circular path of radius r in a magnetic field B . The entire process takes place in a sealed glass tube in which the path of the electrons can be directly observed. During its manufacture, the tube was evacuated, and a small amount of mercury was introduced before the tube was sealed off. As a result, there is mercury vapor in the tube. When electrons in the beam have sufficiently high kinetic energy (10.4 eV or more), a small fraction of them will collide with and ionize mercury atoms in the vapor. Recombination of the mercury ions, accompanied by the emission of characteristic blue light, then occurs very near the point where the ionization took place. As a result, the path of the electron beam is visible to the naked eye.

The tube is set up so that the beam of electrons travels perpendicular to a uniform magnetic field B . B is proportional to the current I through a pair of large diameter coils (so-called “Helmholtz Coils”) in which the coil separation is selected to produce optimum field uniformity near the center.

2 Experimental Apparatus

2.1 The Sealed Glass Tube

Figure 1 shows the filament surrounded by a small cylindrical plate. The filament is heated by passing a current directly through it. A variable positive potential difference of up to 100 volts is applied between the plate and the filament in order to accelerate the electrons emitted from the filament. Some of the accelerated electrons come out as a narrow beam through a slit in the side of the cylinder. The entire tube is located inside a set of coils, which produce a uniform magnetic field B perpendicular to the electron beam. The magnitude of the field can be adjusted until the resultant circular

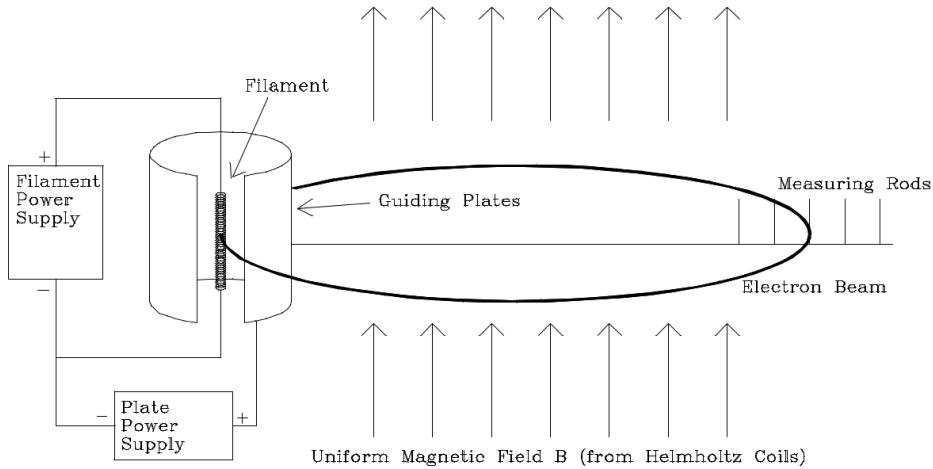


Figure 1: Diagram of the Interior of the Sealed Glass Tube

path of the electron beam just reaches one of the measuring rods. These rods are located along a cross bar, which extends from the cylinder in a direction perpendicular to that in which the electron beam was emitted—i.e., along a diameter of the circular orbits.

2.2 The Helmholtz Coils and Uniform Magnetic Field

The magnetic field produced at the position of the electron beam by a current I flowing through the coils must be computed. For a single turn of wire of radius R , the field on the axis at a distance x from the plane of the loop is given by:

$$B' = \frac{\mu_0 R^2 I}{2(R^2 + x^2)^{3/2}}. \quad (1)$$

For the arrangement in Figure 2, there are two loops with N turns each, separated by a distance equal to the coil radius R . The coils contribute equally to the field at the center:

$$B_I = \frac{\mu_0 R^2 N I}{[R^2 + (R/2)^2]^{3/2}} = \frac{4\pi \times 10^{-7} N I}{R(1 + \frac{1}{4})^{3/2}} = \text{constant} \times I \text{ Tesla} \quad (2)$$

where $N = 72$ is the number of turns of each coil and $R = 33 \text{ cm}$ is the radius of the coils used.

This arrangement, called a pair of **Helmholtz coils**, yields a remarkably uniform field in the region at the center.

The net field B in which the electrons move is not B_I alone, but the resultant of the earth's field B_e and B_I . If the equipment is oriented so that the field of the Helmholtz

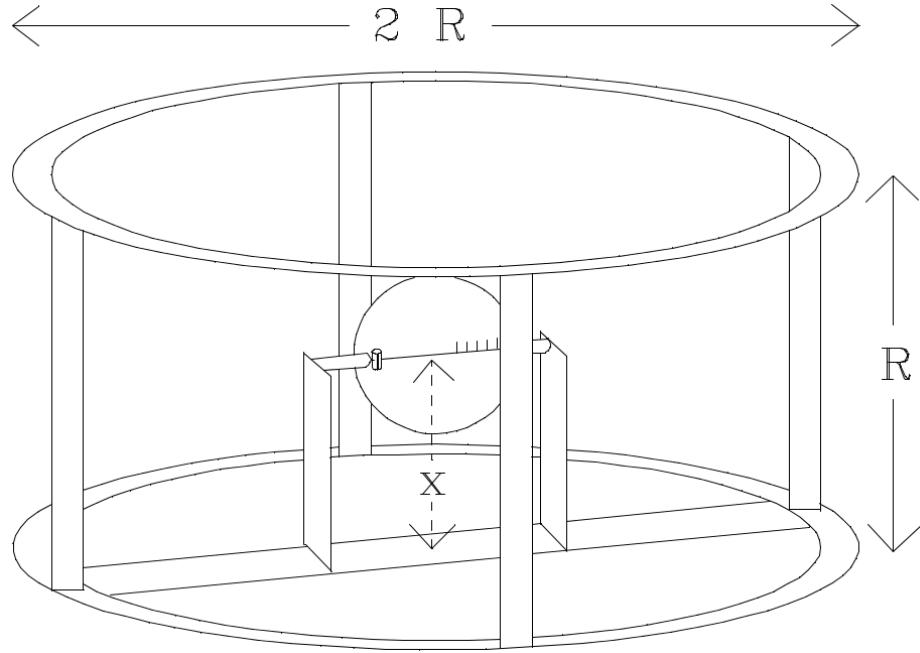


Figure 2: Helmholtz coils used to produce a uniform magnetic field.

coils is parallel to that of the earth, and if the current through the coils causes B_I to be directed opposite to B_e , then

$$B = B_I - B_e \quad (3)$$

2.3 The Trajectories of the Electrons in the Glass Tube

If an electron of charge e and mass m starts nearly from rest and is accelerated through a potential difference V to a final velocity v , then

$$\frac{1}{2}mv^2 = eV \quad \text{or} \quad \frac{e}{m} = \frac{v^2}{2V} \quad (4)$$

If the electron then enters a uniform magnetic field B which is perpendicular to its velocity, it will move in a circular orbit of radius r , where

$$\frac{mv^2}{r} = evB \quad \text{or} \quad \frac{e}{m} = \frac{v}{Br} \quad (5)$$

If it were possible to measure the velocity directly, then e/m could be determined by measurements of either the electric or magnetic field alone. Since a direct measurement of v is not feasible in this experiment, e/m can be determined from the combination of electric and magnetic fields used. Specifically, by eliminating v from equations (4) and (5), e/m can be expressed directly in terms of V , B , and r .

Instead of determining e/m from a single measurement of r for given values of V and B , however, it is preferable to measure the variation of r with B (or I) at fixed values of V . In particular, the data can be presented in convenient form by plotting the **curvature** $1/r$ as a linear function of I .

$$\frac{1}{r} = \sqrt{\frac{e}{m}} \frac{1}{2V} B_I - \sqrt{\frac{e}{m}} \frac{1}{2V} B_e \quad (6)$$

Derive (6) from (3), (4), and (5) for your report before coming to laboratory.

Note that equations (4), (5), and (6) apply strictly only to electrons with trajectories on the *outside edge* of the beam – i.e., the most energetic electrons. There are two reasons why some electrons in the beam will have less energy:

1. There is a small potential difference across the filament caused by the heating current. Only electrons leaving the negative end of the filament are accelerated through the whole potential difference V .
2. Some of the electrons in the beam will lose energy through collisions with mercury atoms.

3 Procedure

3.1 Orientation of the Coil and Tube Setup

For reasons already explained, we would like to orient the Helmholtz coils such that their axes are parallel to the ambient magnetic field.

Please exercise extreme care in the following section as you align the Helmholtz coils: the cathode ray tube is very delicate and may break if the coil support is not very firmly secured and falls. It is recommended that two people handle the coil frame at all times while the coils are being aligned. Also take care not to touch the dip needle, which is easily bent.

In order to align the coils so that their axis is aligned with the ambient magnetic field proceed as follows:

- With the coils in the horizontal position, rotate the horizontal arm of the dip needle support so that the needle itself and the plastic protractor are horizontal. Avoid touching the needle or plastic protractor, and instead turn the arm using the attached lever.
- Allow the compass needle to come to a rest. It is now pointing in the direction of the horizontal component of the ambient field.

- Rotate the entire frame by turning the base, until the compass needle is aligned along the 90° - 270° line on the protractor. The horizontal axis of the cathode ray tube (coaxial with the metal rod) is now aligned with the horizontal component of the ambient field.
- Rotate the horizontal arm of the dip needle support so that the needle itself and the plastic protractor are now in a vertical plane. Avoid touching the needle or plastic protractor, and instead turn the arm using the attached lever.
- Allow the needle to come to a rest. It is now pointing in the direction of the ambient field.
- Loosen the wingnut and **gently** raise one side of the coils. You want to increase the angle until the dip needle is aligned along the 0° - 180° line on the protractor. **Hold the wooden frame rather than the coils as you raise the setup.**
- Securely **tighten the wingnut so that the coils remain in position**. One person should be supporting the frame while a second person tightens the wingnut.
- The coil axis should now be aligned (coaxial) with the ambient magnetic field.

3.2 Preliminary Adjustments

The supplies and controls for the Helmholtz coils and the filament are permanently wired on a board and are designed to minimize the possibility of damage to the tube or coils. Locate each control, and note the qualitative effects observed when the control is varied. In particular:

- (a) Figure 1 shows that the filament and its associated lead wires form a small loop. Since a 4 amp current is required to heat the filament, this loop creates a measurable field. The filament coil reversing switch permits you to study the effect of this field. The effect can be minimized in the experiment by rotating the tube slightly in its mounting so that the electrons come out parallel to the plane of the coils.
- (b) Note the direction of the coil current for each position of its reversing switch by using the dip needle to check the direction of the resultant field. Knowing the field direction, check the sign of the charge of the particles in the deflected beam. Also, determine whether the earth's field adds to or subtracts from the coil field.
- (c) The beam will have a slight curvature in the earth's field when the coil current is zero. Make a preliminary measurement of the earth's field by adjusting the coil current to remove this curvature. The special Meter Switch and low current meter (200 mA) will enable you to measure the relatively small current needed, and the straight line trajectory can be checked by comparison with the light emitted from the filament.

3.3 Measurement of the Circular Orbits

With the accelerating voltage at an intermediate value, the current in the Helmholtz coils can be adjusted so that the outside edge of the beam strikes the outside edge of each bar in turn. Measure field current as a function of radius for the highest voltage V which allows you to adjust the beam with respect to all five bars.

For one measurement, test the reproducibility of the current setting as an aid to error analysis. The tube manufacturer supplies the following values for the *diameters* from the filament to the *outside* of each bar in succession:

$$6.48 \text{ cm}, \quad 7.75 \text{ cm}, \quad 9.02 \text{ cm}, \quad 10.30 \text{ cm}, \quad 11.54 \text{ cm}$$

3.4 Calculations

Plot a graph of $1/r$ versus I , and draw the straight line that gives a best fit to the five measured points. Use equation (6) to calculate e/m from the slope of this line. Write your report up to this point and then proceed to the following:

3.5 Further Considerations

- (a) Calculate the percentage difference between your value of e/m and the accepted value which is 1.758×10^{11} coulombs/kg. What do you think is your largest source of error? Can you account for this much error by a numerical estimate?
- (b) Calculate the actual maximum velocity of the electrons in your beam.
- (c) Make another run with a lower accelerating voltage. Do you find a consistent value of e/m ?
- (d) Compare the intercept of your graph with the value of coil current you obtained by balancing the earth's field. If these numbers are not roughly the same, you may have made an error. Note that this current is not an important number, but it makes a good check on your technique. Physicists frequently check the consistency of their data by computing numbers which they do not "need".
- (e) Calculate the actual value of B_e .
- (f) Test the maximum error in your readings which could be caused by a change in the field of the nearest neighboring coil.

Experiment 2-7

Polarization and Interference

1 Introduction

You have already investigated some wave properties in last semester's Experiment 9. This lab will address three concepts of waves: polarization, interference, and diffraction. We will conduct this experiment with light, but the general concepts apply to many other waves as well.

2 Theory

2.1 Polarization

Polarization occurs only in transverse waves¹. By definition, the oscillation is perpendicular to the direction of wave propagation in a transverse wave. Since in three-dimensional space there are three independent dimensions, there are two possible perpendicular directions needed to describe any oscillatory motion perpendicular to the direction of motion. If you think about a transverse wave in a string, the oscillation could be up-and-down or left-and-right, while the wave front propagates forward along the string. The two independent and perpendicular directions are called the two polarization states. (Or in coordinates, if a wave propagates along the z -axis, there are two polarizations perpendicular to the wave motion—one along the x -axis and one along the y -axis.) Polarization cannot occur in longitudinal waves because the direction of oscillation has only one possibility, along the direction of motion.

You may argue that you can make oscillations happen on a string in more than just the two directions described above—and you are right! For example, you can make oscillations diagonally from the upper left to the lower right or you can even make circular oscillations by moving your hand continuously in a circle, producing a corkscrew pattern along the string². But no matter which wave you make, the oscillations can always be represented as two superimposed waves with displacements along the two mutually perpendicular polarization axes. Any vector can always be described by decomposing it into its x - and y - components. Here, the relevant vector is the oscillating displacement from the z -axis. The amplitudes of the resolved x and y components indicate how much of the wave is in each polarization state.

¹Examples of transverse waves are oscillations on a string, water waves, or the electromagnetic radiation we treat here. Sometimes you may also hear the term polarization used in the context of transverse- and longitudinal polarization.

²This wave is sometimes called left- or right- circular polarization.

For electromagnetic radiation, the oscillating quantities are the electric (and magnetic) fields. These undulate perpendicular to the direction of the wave. Like waves in a string, electromagnetic waves are transverse waves and have polarization properties. The polarization state describes the axis along which the electric field in the wave oscillates. Visible light is electromagnetic radiation in a particular frequency-wavelength regime, with large frequency of oscillation ($\sim 10^{15}$ Hz) and small wavelength ($\sim 0.5 \times 10^{-6}$ m).

2.2 Polarizer

Usually light is produced (e.g. incandescent light from a candle or a light bulb) by hot electrons and atoms, which are moving or are oriented in random directions, so that the wave is **unpolarized**. This means that the electric field associated with the wave oscillates in random directions (though always at right angles to the direction of the propagation of the light). So there is no favored direction for polarization. It is as if you were moving your hand in random directions at the end of a string. Note that this is very different from the cases of waves in a string described above. Even in circular polarization the oscillations were not random.

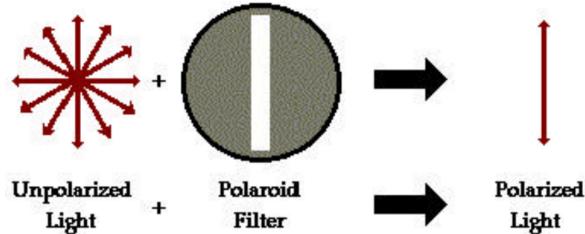


Figure 1: Polarization Filter and Creation of Polarized Light

A polarization filter, shown in Figure 1, is a device with a specified direction, called the polarization axis. The filter absorbs all the wave's incoming electric field perpendicular to the polarization axis and transmits the electric field parallel to the polarization axis. So no matter what the polarization of the light shining into a polarizer might be, the light coming out of the polarizer is always polarized in the direction specified by the polarization axis³. An unpolarized wave just prior to the polarizer has an electric field (**E**) pointing randomly in all directions perpendicular to the wave direction, as shown in Figure 2. The polarized wave after the polarizer has only half the energy content, and the **E**-field is oscillating solely along the polarization axis. This electric field has the same value as the component along that axis just before the polarizer.

³Such light is called linearly polarized light.

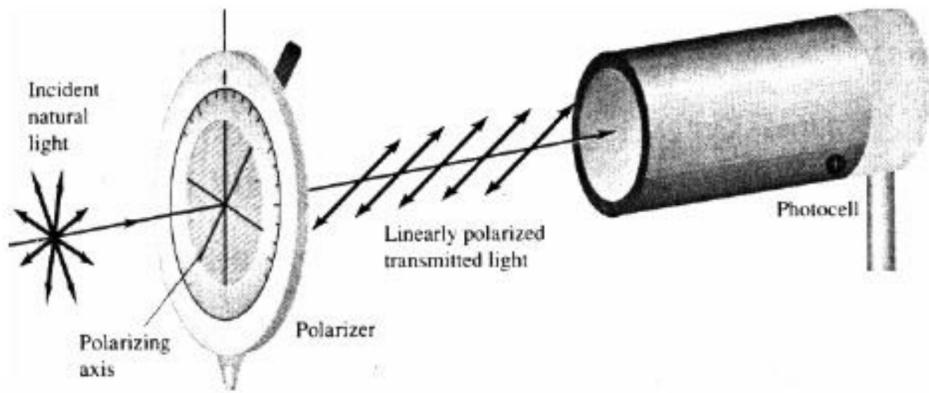


Figure 2: Detection of Polarized Light

2.3 Interference—The Double Slit

As you've learned in class, Young's double slit experiment is a demonstration of interference and diffraction. A double slit, as illustrated in Figure 3, permits what remains of an incoming wave (from the left) to travel to a distant screen (to the right) along two different paths with different lengths. The light waves from the two slits interfere, resulting in an interference pattern of bright (constructive) and dim (destructive) patches as viewed on the screen. Let's see how the double slit works using Figure 3.

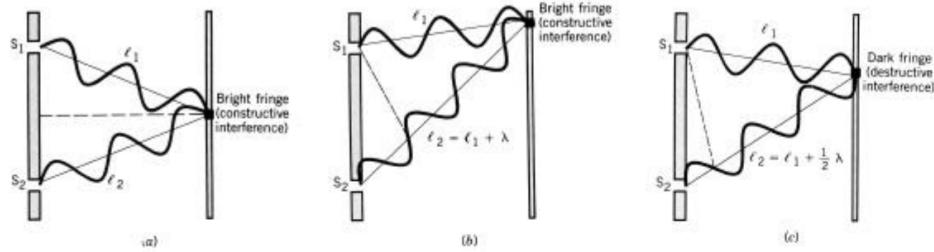


Figure 3: Young's Double Slit Experiment

We shine light from the left onto the double slit, which allows two light waves to propagate from the two slits. Straight-ahead they will always be in phase since they travel the same distance to the screen. But when the two waves propagate at an angle θ , they cover different distances to reach a specific point on the screen. (We assume that the two rays are parallel, a good approximation since the screen is effectively infinitely far away compared to the distance between the two slits.)

As you can convince yourself by applying the rules of trigonometry to the triangle in Figure 4, the difference in path (Δl) between the two rays is, for slit separation, d ,

$$\Delta l = d \sin \theta \quad (1)$$

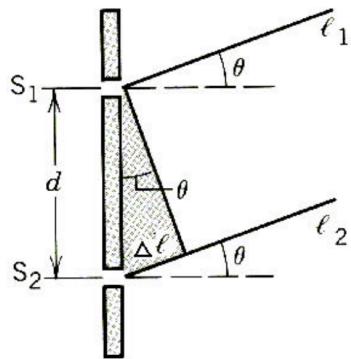


Figure 4: Geometry of the Double Slit

So what quantitatively determines a point of maximum intensity on the screen? Maximum intensity occurs if the two waves are precisely in phase at that point on the screen. This follows from what you learned about constructive and destructive interference in the lab on standing waves. As you saw, two waves can add up constructively or destructively. In the case of maximal constructive interference, the waves combine at the screen so that a crest on one wave coincides with a crest on the other wave and a trough coincides with a trough. In the case of maximal destructive interference, the waves combine so that a crest on one wave coincides with a trough on the other wave.

The waves will be in phase when they meet at the screen if the path difference between the two waves is an exact number of wavelengths, m , where m is an integer called the order of the maximum. This requirement, expressed mathematically, is

$$\Delta l = m\lambda \quad (2)$$

Note that if $m = 0$, there is no difference in the path lengths, resulting in a maximum straight ahead, $\theta = 0$, as shown in Figure 5. Subsequent maxima occur for subsequent positive and negative integers, $m = \pm 1, \pm 2, \pm 3, \dots$

Combining the two equations for Δl , we get for the specific angles, θ_m , that give maximal constructive interference:

$$d \sin \theta_m = m\lambda \quad (3)$$

If we know the order m of a particular maximum, we can determine the wavelength of the light incident on the double slit by measuring the angle of the maximum on the screen. Conversely, we can determine the order of a particular maximum, m , from the wavelength of the light and the angle.

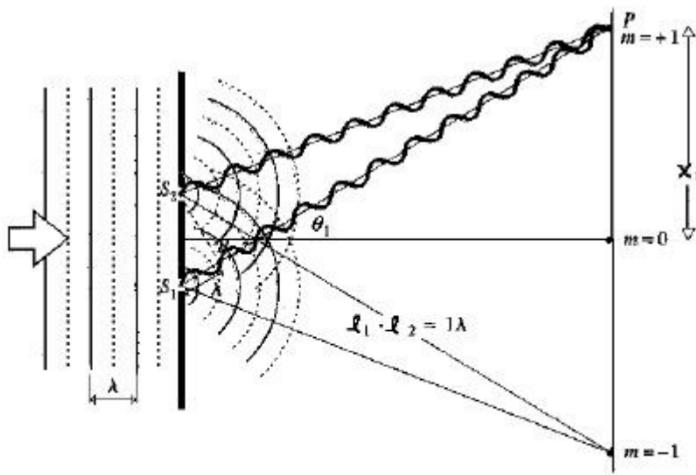


Figure 5: Interference of Light Coming from Two Small Slits

2.4 Intensity of the Double-Slit Maxima

You will soon note that the intensities for subsequent maxima are different. What causes this? Just as the double slit creates an interference pattern, each single slit creates a diffraction pattern due to interference from waves coming from different parts of a single slit as described in section 39-6 of the lecture course text⁴. This occurs because each individual slit has a finite width. The overall (more realistic) pattern shown in figure 6 is a superposition of the single slit and double slit interference patterns. The narrow spikes (of interest for this experiment) are due to the double slit interference and the envelope, or large-scale pattern, is due to the single slit diffraction caused by each slit.

2.5 Light from a Laser

In the double slit experiment we use a laser as the light source. Laser light differs from incandescent or fluorescent light in two important ways. First, lasers emit intense light of a single wavelength or frequency (monochromatic). Second, laser light is coherent, meaning that all light waves are in phase.

Let's illustrate the differences between coherent and incoherent light through a simple familiar example. Contrast a crowd of people cramming the city streets during rush hour with a regiment of soldiers at a parade. In the midtown crowd, even if everybody were to step at the same frequency (which they don't!), each pedestrian's steps are independent of what others are doing. But soldiers at a parade move in step, or in phase, with a similar frequency, taking their steps at exactly the same time. Laser

⁴Sears, Zemansky, and Young, **College Physics 7th edition**, Addison-Wesley, p. 903-906.

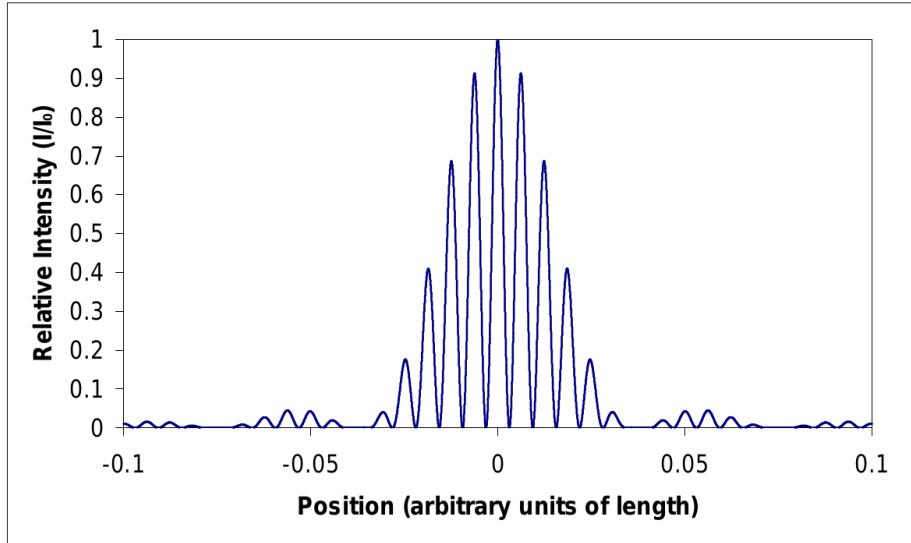


Figure 6: Intensity Pattern of the Double Slit. The Many Peaks of the Narrow Spikes Correspond to the Maxima in $d \sin \theta_m = m\lambda$

light is like the (coherent) soldiers at a parade, whereas light from a light bulb is like the (incoherent) civilian crowd at Times Square.

Although incandescent light sources⁵ can produce diffraction and interference patterns, a laser is better suited to illustrate coherent phenomena.

3 Experiments

3.1 Polarization

Suppose a wave has linearly polarized light characterized by an electric field vector, \mathbf{E}_0 , at an angle, ϕ , relative to the direction of a polarizer. The electric field that is transmitted through the polarizer is $\mathbf{E} = \mathbf{E}_0 \cos \phi$, the component that was not absorbed. Since the energy intensity in the wave is proportional to \mathbf{E}^2 , then the intensity passing through two polarizers will be proportional to the square of the cosine of the angle between the two polarizing axes.⁶

$$I = I_{\max} \cos^2 \phi \quad (4)$$

In this experiment, we use an incandescent unpolarized light source and then polarize it by first passing it through a polarizer. The electric field vector \mathbf{E}_0 will then

⁵This is easily visible in the colors produced by soap bubbles.

⁶This relation is sometimes referred to as Malus' Law, after Etienne Malus, who also discovered the polarization of reflected light in the early nineteenth century.

point along the same direction as the polarizing axis of the first polarizer. As the now polarized light falls on the second polarizer, only the component of \mathbf{E}_0 along the polarizing axis of the second polarizer will make it through unaffected.

As figure 7 illustrates, an electric field strength of only $\mathbf{E}_0 \cos \phi$ will make it through the second polarizer, so that the new intensity will be given by Equation 4 above. A human eye or a photometer can detect this change in light intensity.

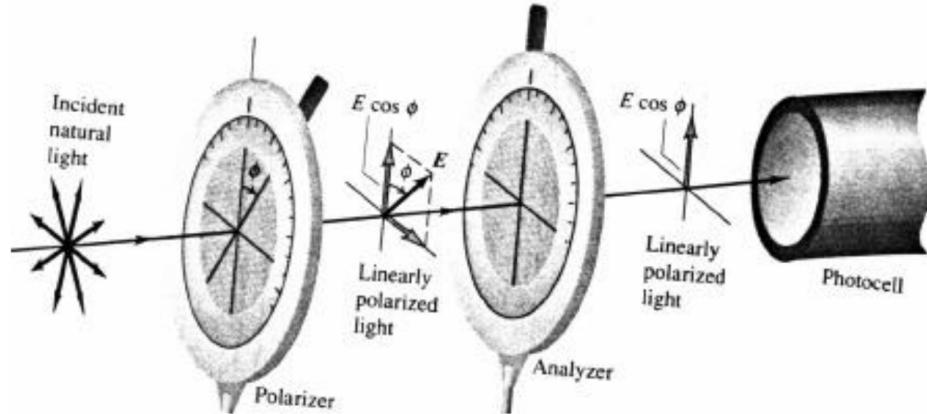


Figure 7: Creation and Detection of Polarized Light

3.2 Double Slit

In this part of the experiment we use a laser instead of incandescent light source. Laser light shines onto a double slit, which produces an interference pattern of bright and dark lines on the wall or a piece of paper beyond the double slit.

We know that the angle, θ_m , at which the m^{th} order bright spot appears is given by

$$d \sin \theta_m = m\lambda \quad (5)$$

For small θ_m (in radians), the approximation $\sin \theta \approx \tan \theta \approx \theta$ is valid⁷. So on a screen far away from the double slit, we observe maxima located close to the optical axis (where the 0^{th} order maximum hits the screen) approximately given by

$$\sin \theta_m \approx \tan \theta_m = \frac{x_m}{L} \quad (6)$$

where x_m is the distance between the m^{th} order maximum and the 0^{th} order maximum, and L is the distance between the double slit and the screen. We then find

$$\frac{x_m}{L} = m \frac{\lambda}{d} \quad (7)$$

⁷See for yourself by punching $\sin 0.1$ and $\tan 0.1$ into your calculator (using radiant measure)!

or for adjacent maxima

$$\frac{\Delta x}{L} = \frac{\lambda}{d} \quad (8)$$

with Δx equal to the distance between the maxima. Since Δx and L are easy to measure, we can straightforwardly determine d (or λ) given λ (or d).

4 Specifics of the Experiment

SAFETY NOTE:

*Although the laser is of low intensity for most situations, it can be dangerous in certain circumstances unless used carefully. In particular, **do not** use the laser in such a way that it can shine into any person's eye (The warning label on the laser states "Do not stare into Beam!"). When you are not actually using the laser, turn it off or close the beam shutter at the front of the laser.*

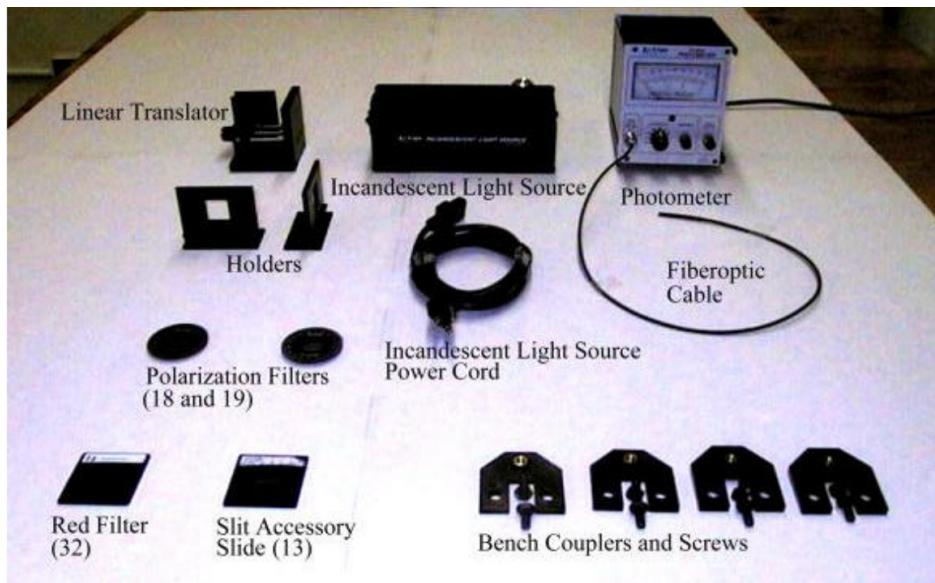


Figure 8: Equipment Components

5 Equipment

Linear Translator: A holder for the fiberoptic cable that can be moved transversely in small increments, allowing a positioning with better than 0.1 mm precision.

Fiberoptic Cable: Cable made from glass, quartz, plastic or a similar material, which carries light waves. The cable uses total internal reflection at the walls to provide very high efficiency (or low losses) between the cable ends.

Photometer: Device to measure the total power contained in light impinging on it. It uses the photoelectric effect to convert the photons to electrons and various electron amplification techniques to produce a measurable electric current.

5.1 Polarization

- Make sure only the dim incandescent ceiling lights in the room are on.
- The linear translator should be placed at the far end of the track opposite the laser. Make sure that the linear translator is in the middle position (at 2.5 cm).
- Place the fiber-optic cable in the hole of the linear translator.
- Place the incandescent light source about 10-20 cm from the linear translator such that the light shines towards the linear translator.
- Place the polarizers (numbers 18 and 19) in the holders.
- Rotate the polarizers such that the white mark on the holder aligns with the 0° angle reading on the polarizer.
- Place the two polarizers between the incandescent light source and the linear translator.
- Put the photometer to the lowest sensitivity (the 1000 setting) and use the zero adjust knob to make sure that the pointer is at the 0 position when the incandescent light source is turned off.
- Turn the incandescent light source back on and adjust the sensitivity of the photometer such that the needle is at the highest position without being at the maximum reading. The value of the sensitivity corresponds to the maximum value on the analog scale.
- Record this measurement as the intensity for 0° difference between the two polarizing axes.
- Measure the intensity for angles between 0° and 90° in increments of 10 degrees. If the readings become too small you may have to switch the photometer to a setting with a higher sensitivity.
- Plot the relative intensity (I divided by I_0 , the intensity at 0°) vs. the \cos^2 of the angle between the axes of the polarizers.
- What is the general shape of your graph? What does this indicate? Did your line or curve pass through the origin?

- What reading do you obtain when the polarizers are at an angle of 90° relative to each other? (This is called the “noise” of your measuring device.) What reading would you expect if there was no noise?
- What is the signal to noise ratio? (The “signal” is the reading at 0° .) What does this tell you?
- If you had a lot of background light, how could you reduce the influence it would have on your results (by changing the setting or by changing your data analysis)?
- If you had aligned the 90° mark instead of the 0° mark of both polarizers with the white mark on the holder would your results have been any different? What relative intensity would you expect if the angle between the two polarizers was 180° ? 270° ?
- What are the main sources of error? Which do you think contributes most?

5.2 Double Slit

- Before you start this section, be sure you have read the laser safety note above.
- Your TA will demonstrate the interference pattern by projecting it on the wall. Note down your observations!
- Remove the incandescent light source, the polarizers, and the holders from the track, and the fiberoptic cable from the linear translator.
- The laser should already be mounted on the laser holder at the far end of the track. Rotate the dial on the linear translator so that it is in the middle position.
- Turn on the laser. It should propagate through the center of the hole in the linear translator, producing a bright red dot on the wall. If the laser is not properly aligned, alter the position of the laser by adjusting *only the back legs* on the laser holder. Make sure you avoid lifting the laser and shining it in your classmates’ eyes.
- Place the red filter (number 32) onto the linear translator. Carefully slide the fiberoptic cable back into the hole in the linear translator so that the front of the cable sits right against the red filter.
- Place the slide with the slits (number 13) in a holder and place the holder right in front of the laser.
- Adjust the slide position in the holder so that the laser propagates through slit pattern A.

- Note the interference pattern on the front of the linear translator.
- Adjust the photometer sensitivity to an appropriate setting.
- Turn the knob on the linear translator. This should cause the intensity reading on the photometer to change.
- Record the position and intensity for each maximum until the linear translator has moved as far as possible, starting from the 0 position of the linear translator. As you turn the knob on the linear translator take care not to shift the position of the linear translator on the bench.
- Measure and record the distance L from the slits to the tip of the fiberoptic cable. Where is the 0th order maximum located? Why might it be located somewhere other than the middle position of the linear translator?
- Plot the position of the maxima x_m vs. the order m . Determine the slope of your graph and use it to calculate the wavelength of the laser light. The slit separation for slit A is 0.250 mm.
- Construct a graph of the relative intensity, I/I_0 , vs. x_m . Explain the shape of your graph.
- Would the maxima be more spread out or closer together if the slit separation decreased?
- Remove the slit slide and its holder from the track and the fiberoptic cable from the linear translator. Rotate the dial on the linear translator so that it is in the middle position.
- Place the unknown slit apparatus on the track approximately 25 cm from the front of the laser so that the laser propagates through the slit pattern marked with a pink dot. You should see a horizontal interference pattern on the linear translator; the light in the center of the interference pattern should propagate through the hole in the linear translator. Note: If the laser does not propagate through the appropriate slit, slide the unknown slit apparatus towards and away from the laser until the laser strikes the slit pattern at the appropriate place. If the light from the pattern does not propagate through the center of the linear translator, alter the position of the laser by adjusting *only the back legs* on the laser holder.
- Again measure the distance between consecutive maxima as above. Also, measure the distance L from the slits to the tip of the fiberoptic cable. Use this data and the wavelength of the laser you determined previously to determine the slit separation.

- What is the purpose of using the red filter?
- Give the main sources of error.

6 Applications

Some chemical molecules have a property called chirality. This means that the mirror image of the molecule is not identical to the original. For example, no matter how good you might be at 3D puzzles, you cannot arrange images of the molecule and of its mirror image to superimpose precisely. A famous example is the DNA double helix, a highly complex molecule, which carries all genetic information. The molecule always twists in one direction, while its mirror image always twists in the opposite direction. Many such chiral molecules are optically active, and rotate the polarization-axis of transmitted or reflected polarized light⁸. If you shine polarized light into such a sample, you will see that the light leaving the sample has its polarization axis rotated a few degrees clockwise or counterclockwise, depending on the substance.

This effect can be used to determine the concentration of chiral molecules in a solution. For example, you can determine the glucose concentration in a blood specimen in a non-chemical way. The calculation of the concentration is particularly simple, since the change in angle of the polarization axis depends only on the concentration of the chiral substance in the solution, the distance the light travels through the sample, and a constant specific for the substance. Calculating the glucose concentration, therefore, requires only a measurement of the rotation in polarization vector and known constants – a simple task for a rudimentary computer.

⁸For instance, the glucose dextrose turns the polarization axis to the right: “dexeter” is Latin for “right”.

7 Lab Preparation Examples

Polarization:

1. You look at the sky through a polarizer and as you turn the polarizer you see that the sky appears darker or lighter, depending on the position of the polarizer. What does this observation tell you?
2. You look at a light source through a polarizer. As you turn the polarizer you see no change in intensity. Is the light emitted by this source polarized?
3. Unpolarized light passes through two polarizers whose polarization axes are aligned. You note down the observed intensity as I_0 . Now you turn the polarization axis of one of the polarizers by 60° . What intensity do you observe now?
4. Now you add, after the second polarizer, a third polarizer. The third has a polarization axis of 30° relative to the second polarizer. What intensity do you observe now?
5. Does the result depend on which of the two polarizers you rotate in question 3?
6. Would your answer to question 5 change if the light source produced polarized light? Explain!
7. You have two polarizers with an angle of 90° between their polarization axes. Adding a third polarizer between these two polarizers will typically produce light after the final polarizer. At what angle (relative to the first polarizer's axis) would you observe the maximum intensity after the third polarizer?

The Double-Slit:

While performing a double slit experiment, with $L = 1.00 \pm 0.01$ m, $d = 0.5 \pm 0.1$ mm, you obtain the following locations for maxima:

order	-4	-3	-2	-1	0	1	2	3	4
Position (mm)	-4.1	-3.2	-2.0	-0.9	0.0	1.3	2.1	3.2	3.9

8. Calculate x for all of the consecutive maxima and use the $2/3$ estimate to determine the uncertainty of x . Now determine the wavelength of the light (including uncertainty).
9. Using the data above, determine x from a plot of position vs. order. Given the wavelength as $\lambda = 600$ nm and $L = 1.00 \pm 0.01$ m, what is the spacing between the two slits?

8 Appendix for Experts: How a Laser works (Not required!)

Quantum mechanics tells us that light is produced when electrons in an atom make a transition from a higher to a lower state, or energy level. Whenever an electron makes a transition between states, the difference in energy is emitted (or absorbed) as a discrete energy package called a photon⁹, which carries (with other photons) the electromagnetic wave. The transition from a higher to a lower state can happen in two different ways. First, it can happen randomly, i.e. the electron just falls down spontaneously. This is what happens when a candle produces light. Electrons that make such random transitions result in emitted photons which are incoherent, or out of phase with one another.

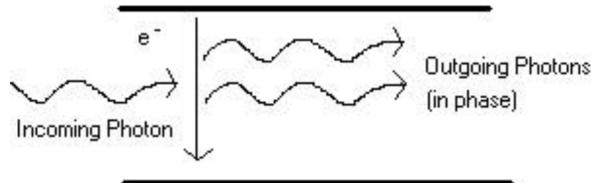


Figure 9: Process of Stimulated Emission

The second way this process can happen is called stimulated emission. If you already have a field of photons that are in phase and have the same frequency, and that frequency corresponds exactly to the transition energy, the emission is no longer random. When an electron falls to the lower energy level, a photon is emitted precisely in phase with the already existing photons, as shown in Figure 9.¹⁰ In our analogy of soldiers marching on parade, new entrants to the parade are required to start walking in step with the others (or they would miss out on the parade).

For stimulated emission to work, the upper atomic state must be full of electrons and the lower one must be almost empty. This is called an inversion, since atoms (at ordinary temperatures) tend to be in their lowest energy state. The inversion is actually achieved by pumping the electrons from the lower state into a third transient state, with an energy even higher than the second laser state, so that the electrons fall

⁹See the experiment on the Photoelectric Effect.

¹⁰Photons are of a particle type called bosons. Such particles are highly “social” and like to be in exactly the same energy state and the same phase with their friends. This contrasts with another particle type, called fermions, which hate being in the same state as their associates (Pauli Exclusion Principle). Electrons are fermions, and this antisocial property accounts for the unique state of each electron in an atom. If electrons were bosons, all electrons in an atom would be in the ground state and we would not have chemistry!

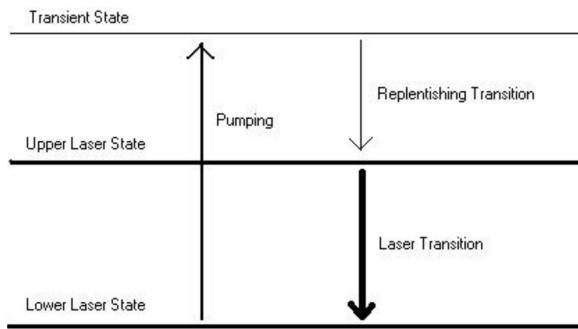


Figure 10: The basic three - level scheme for laser operation.

down to replenish the upper laser state. So real lasers use at least three energy levels; efficient lasers usually use even more¹¹.

Finally, the setup is placed in a cavity with a mirror at one end and a semi-transparent¹² mirror at the other end. This increases the efficiency of the laser by maintaining a high density of laser light within the cavity. The little leakage emitted from the semitransparent mirror is actually what we use as the laser beam.

¹¹Semiconductor lasers (or laser diodes) work somewhat differently, but the underlying concepts of stimulated emission and inversion are the same.

¹²These semi-transparent mirrors are still highly reflective. They reflect about 99% of incident light to keep most of the light inside the cavity.

Experiment 2-8

The Photo-Electric Effect

The photoelectric effect was the first direct experiment to demonstrate the quantum nature of electromagnetic energy transfer to and from matter. The purpose of this experiment is to verify Einstein's theory of the photoelectric effect, to measure Planck's constant, and to determine the work function for a surface.

1 Physical Principles

The photoelectric effect may be understood as a consequence of quantum mechanics, according to which light is carried by discrete bundles of energy (quanta) called photons. Photons have energy hf , where h is a universal constant called Planck's constant and f is the frequency of the classical electromagnetic wave. When a photon is incident on a metallic surface, it interacts with an atom in the metal and transfers all its energy to one of the atom's electrons. This electron may then escape through the electric field at the surface, which keeps less energetic electrons inside the metal. The emerging electron then has energy equal to the energy of the photon minus the energy W lost in escaping the metal. W , the **work function** of the surface, is a material-dependent constant. Since electrons also lose energy in collisions with other electrons before emerging, we may only specify the maximum possible energy for an electron liberated by light of frequency f from a metal. If the material work function is W , this maximum energy¹ is

$$E_{\max} = hf - W \quad (1)$$

According to quantum mechanics, a plot of E_{\max} against f should give a straight line with slope h , as shown in the figure 1. The line should intercept the f -axis at $f_0 = W/h$. This value f_0 , the lowest frequency that can eject an electron, is called the **threshold frequency**.

2 Historical Perspective

Early attempts to describe the photoelectric effect used classical electromagnetism. These studies concluded that the stopping voltage, or potential needed to keep the highest energy electrons from reaching the electrode, should depend only on the intensity of the light but not on the frequency. (In other words, the electrons should be ejected with more kinetic energy if you shine a bright light rather than a dim light.)

¹The maximum electron energy is measured in photoelectric experiments by finding the potential (V , in Volts) at which a decelerating electric field causes all electrons (charge e) to stop before reaching the electrode. Therefore, $E_{\max} = eV$.

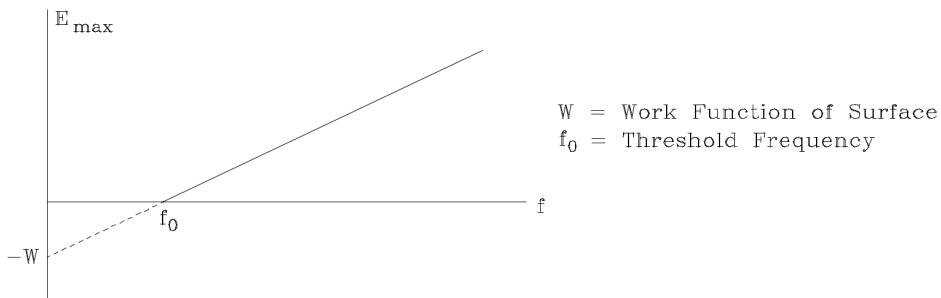


Figure 1: Maximum kinetic energy of the ejected electrons as a function of the frequency of the incident electromagnetic radiation

But it turns out that experimental evidence contradicts these predictions! If you shine a light with any frequency below f_0 , no electrons will be set free no matter how intense the light.

Albert Einstein was the first (in 1905) to analyze the photo-effect problem based on a concept of light as quantized energy packages. The idea earned him his Physics Nobel-Prize in 1921. The idea, though simple, was revolutionary at the time. Einstein hypothesized that light consists of little discreet energy packages, or quanta, which behave like particles called photons. An individual photon cannot be divided, but it can be totally absorbed under appropriate circumstances. The size of the photon energy quantum is just determined by the frequency f of the electromagnetic light, as $E_{\text{photon}} = hf$. For the electron that absorbs a photon to leave the metal requires a minimum energy W to set it free. (The value of W depends on the metal.) So after the electron leaves the metal, it still has energy $hf - W$, which it carries as kinetic energy. This is where equation (1) comes from.

Does this mean that all the classical electromagnetism you learned so far this semester is wrong? No! Electromagnetism is still a good theory that makes quantitative and correct predictions about the real world, so long as you deal with many photons and many atoms (i.e. in the macroscopic world). Only when it comes to the behavior of a single or few atoms, photons and/or electrons does one need to use Quantum Mechanics.²

3 Experimental Apparatus

The photoelectric effect in this experiment occurs at the cathode of an IP39 phototube. The phototube is constructed as shown in Figure 2. The IP39 phototube is designed to

²We said that classical electromagnetism claims that higher intensity means more energy in the light. This is still true. But classically, more intensity just means more photons, not photons of higher energy.

function as a diode where light, incident on the tube, illuminates the cathode and causes photo-emission of electrons. The central wire (anode) is shielded from illumination.

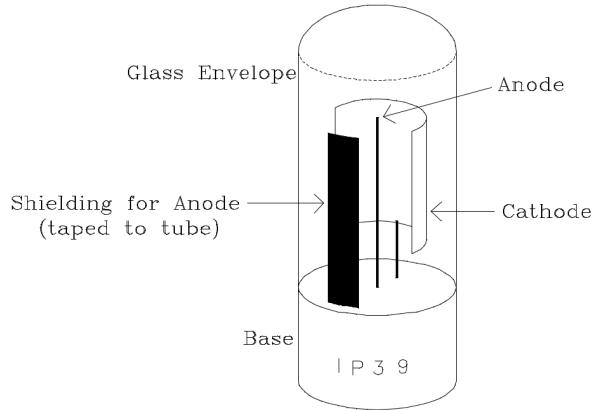


Figure 2: The IP39 Phototube

If the anode is connected through an external circuit to the cathode, a current will flow whenever the cathode is illuminated with light of frequency greater than f_0 . If a potential V is applied so that the anode is *negative* with respect to the cathode, as illustrated in Figure 3, the emitted electrons will lose kinetic energy as they travel through the potential field. Electrons with kinetic energy less than eV will not be able to reach the anode, and the current in the external circuit will be less than it would if $V = 0$. For a potential $V = V_S = E_{\max}/e$, called the **stopping potential**, the current will drop to zero.

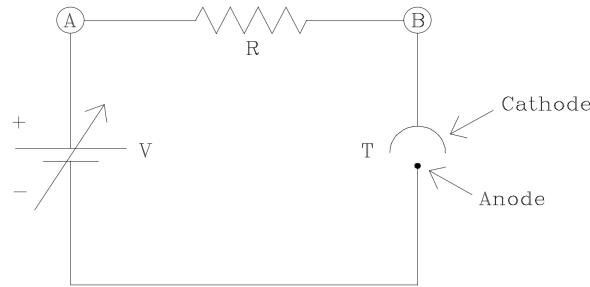


Figure 3: Simplified Phototube Circuit

In more detail, as V is varied, there are three possible cases to consider (see Figure 3):

Case I: $V < V_S$. In this case, the anode (the symbol for the anode is the dot) will collect electrons. The tube T will act like a battery of greater strength than V , and the electrons will flow around the circuit in a clockwise direction. The current will cause an iR drop across R . Thus, A and B will be at different potentials.

Case II: $V = V_S$. At this critical value of V , electrons are just not energetic enough to reach the anode. No current flows, and A and B are at the same potential.

Case III: $V > V_S$. One might think that electrons would flow counterclockwise. This can happen if there is light striking the anode. If no light is permitted to strike the anode, no electrons are emitted from it and thus the tube T acts like an open circuit; A and B are again at the same potential.

In principle, the stopping potential (for a given frequency of incident light) could be measured by finding with a voltmeter the value of the adjustable negative voltage V at which the current passing through an ammeter in series with the phototube goes to zero. However, the number of photoelectrons liberated by available light sources produces currents of less than 2 microamperes, too small to be measured by simple meters. A more sensitive way is needed to determine the presence or absence of current in the phototube circuit.

A Photoelectric Module has been designed using high gain current amplifiers for easy measurement of the stopping potential for the IP39 phototube for several different frequencies of incident light. A schematic diagram of the phototube circuit using the Photo-electric Module is shown in Figure 4. The High Gain Difference Amplifier is used to detect any current flowing through the resistor R and consequently through the phototube. Both the High Gain Difference Amplifier and the 0-3 V meter are designed to draw negligible current from the phototube circuit. If the **Zero Adjust** control has been correctly set, then a reading of zero on the 5-0-5 meter means that no current is flowing through the phototube. The stopping potential applied to the phototube is set using the **Voltage Adjust** control and read with the 0-3 V meter.

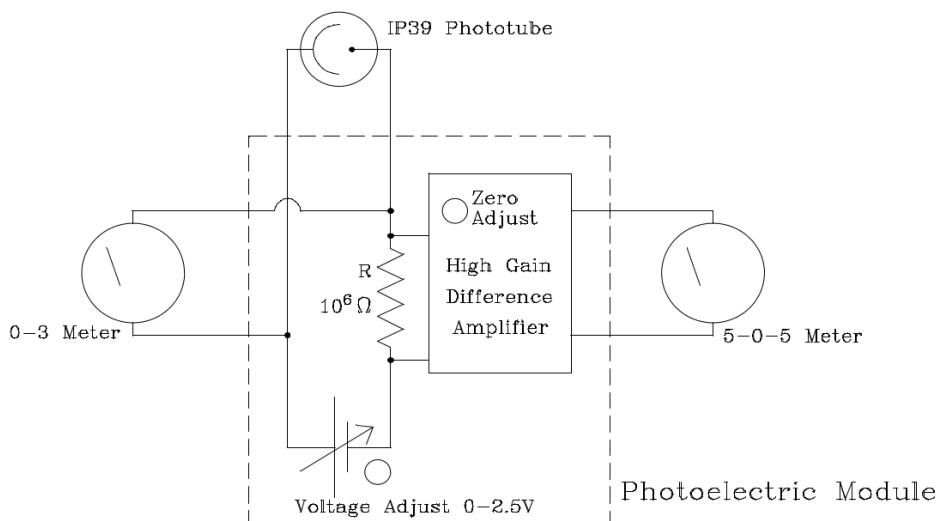


Figure 4: Schematic Diagram of the Phototube Circuit Using the Photo-electric Module

The overall arrangement of the apparatus is shown in Figure 5 below. The Phototube circuitry is housed in the Photoelectric Module, which is connected externally to the two voltmeters as well as the Phototube. The Phototube itself is attached over a window in the housing of a mercury arc source, which provides light that is passed through selected filters and then strikes the cathode of the tube.

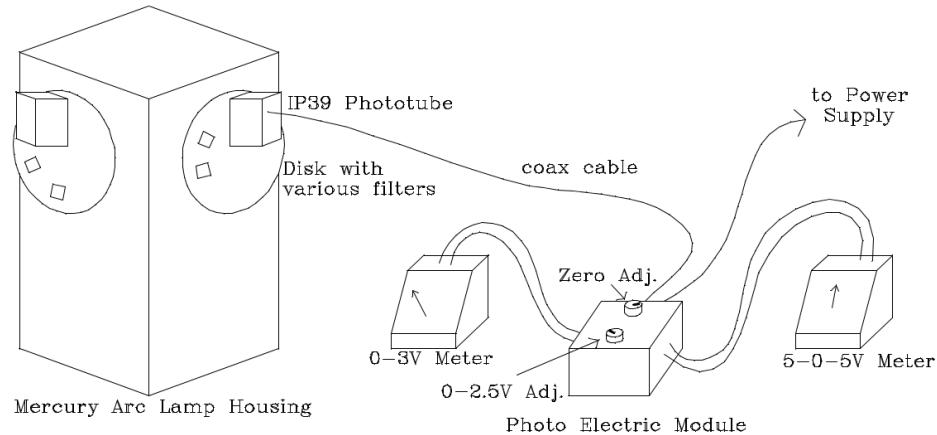


Figure 5: Arrangement of the Experimental Apparatus.

In order to vary f , one would like to use an ideal **monochromator**, i.e., a source which produces light at a well-defined but variable frequency. A discharge tube, such as the mercury arc source, produces a series of discrete frequencies, i.e., **spectral lines**. A prism or grating spectrometer could be used to transmit any one of these lines while rejecting the others, but the narrow slits on such an instrument will not provide sufficient intensity for this experiment. Therefore, we use instead a set of bandpass optical filters, each of which transmits only a small band of frequencies. The arrangement is shown in Figure 6.

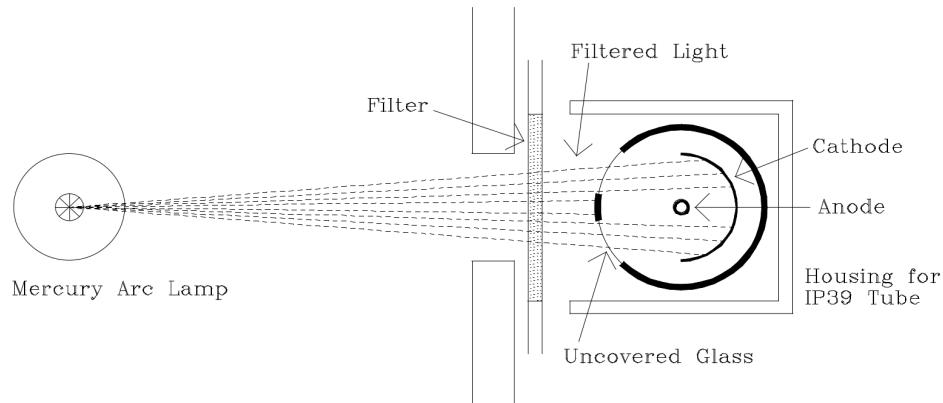


Figure 6: Light source and Phototube Arrangement.

Figure 7 presents a graph indicating the percentage of light transmitted as a function of wavelength for the four filters to be used. On the same graph the principal wavelengths in the visible spectrum of mercury are indicated as straight lines.

Although a filter may transmit more than one line, only the most energetic photons determine E_{\max} . (Note that the vertical scale of Figure 7 is logarithmic; the *small percentage* of the intensity of the most energetic lines passed by two of the filters is not sufficient to affect the results of this experiment.)

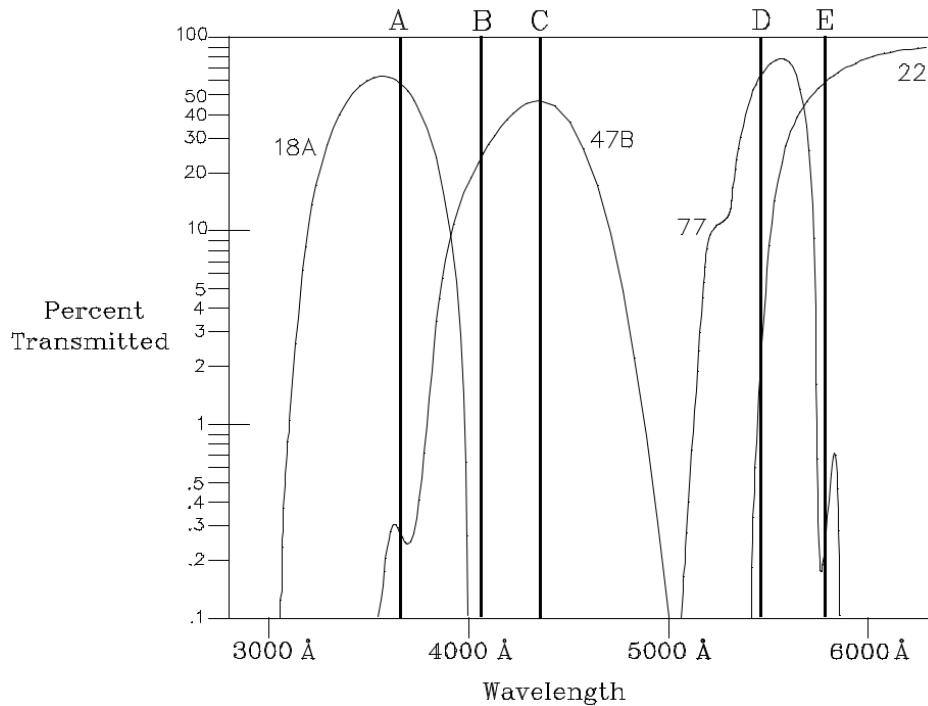


Figure 7: Percentage of Light Transmitted by Four Bandpass Filters.

Mercury Lines Shown on Graph ($1 \text{ \AA} = 10^{-10} \text{ m}$)

- | | |
|-------------------------------------|-----------------------------|
| (A) Triplet: 3650 Å, 3655 Å, 3663 Å | (D) Singlet: 5461 Å |
| (B) Doublet: 4047 Å, 4078 Å | (E) Doublet: 5770 Å, 5791 Å |
| (C) Singlet: 4358 Å | |

4 Procedure

4.1 Determining Planck's Constant h and the Work Function W of the Phototube

E_{\max} can be measured for a few different frequencies of incident radiation using the mercury arc source and filters as described above. The stopping potential for each filter can be determined using the Phototube circuit shown in Figure 4. Turn on the power supply for the Photoelectric Module and the mercury lamp. Since the High Gain Difference Amplifier drifts slowly with time or temperature, it is important to set the Zero Adjust control before you start to make measurements.

First, turn the **Voltage Adjust Control** to the fully counterclockwise position. (The 0-3 meter should read 0.) Then, the **Zero Adjust** is set so that the 5-0-5 meter reads zero when no light is allowed to shine on the cathode. (The phototube is blocked by turning the filter disk to a no-window position. Note that the window which allows you to see light through the filter *is not the same as the window which allows light into the phototube housing.*) The Zero Adjust setting should be checked periodically.

Now, stopping potential measurements can be made. Rotate one of the filters into position between the light source and the phototube. The 5-0-5 meter should show a large deflection. Turn the Voltage Adjust knob clockwise until the reading on the 5-0-5 meter reads zero or near zero. Once the needle on the 5-0-5 meter stops moving toward zero, the stopping potential can be read from the 0-3 meter. Note: It is very important to carefully monitor the 5-0-5 meter when turning up the voltage across the phototube. *The stopping potential is the minimum voltage that stops the current flow through the phototube.* The stopping potential can be measured for other light frequencies by turning the Voltage Adjust to zero, selecting another filter, and repeating the above procedure. Plot V_S versus f and determine h and W .

4.2 Disproving the Classical Theory

The classical theory (Newtonian mechanics and classical electromagnetism) predicts that if the *intensity* of the light is cut in half, then the stopping potential will likewise be cut in half. Use the filter which is half masked with aluminum foil and see whether the stopping potential remains the same as with the equivalent unmasked filter.

Questions:

1. Why is the anode of the phototube shielded from light?
2. If the light on the phototube contains two different wavelengths λ_1 and λ_2 , which value must be used in the formula for the stopping potential? What is the effect

of light of the other wavelength?

3. Does your value for the work function W of the tube have the order of magnitude you expected? Why did you expect (or not expect) the work function to be of this magnitude?
4. What are the experimental features of the photoelectric effect that Newtonian mechanics and classical electromagnetic theory cannot explain?
5. The surface of the photocathode of the IP39 tube is an alloy of cesium and antimony. Impurities (atoms from the anode, for example) can sometimes become deposited on the cathode surface. What effect would such a deposit have? How would it affect your values for the stopping voltages? Decide from your data whether this effect may be important for your tube.

Experiment 2-9

The Spectrum of the Hydrogen Atom

1 Introduction

In this experiment we will observe the discrete light spectrum one observes from a gas discharge lamp. We will see that the spectrum consists of a collection of sharp, single colored lines. We will be able to measure the wavelength of the light emitted quite precisely, usually better than 1 part in a thousand. Therefore it is crucial to make all calculations to 5 significant figures.

2 Theory

2.1 The Spectrum of the Hydrogen Atom

According to the classical theory of electro magnetism, atoms should radiate continuously and the electrons should fall into the nuclei within a short timespan. Obviously this is not the case since we have stable atoms. Therefore the classical description seems to be wrong.

Niels Bohr suggested an alternative theory for atoms, suggesting that the electrons can only exist in certain energy states and can only jump between these discreet states discontinuously. This discontinuity explains why we see a collection of sharp lines in a gas discharge lamp, and not a continuous spectrum as in a usual light bulb.

Bohr was able to derive the following formula for the energy of the energy levels only using natural constants

$$E_n = -\frac{2\pi^2 me^4 k_e^2}{n^2 h^2} = -\frac{me^4}{h^2} \left(\frac{1}{8\varepsilon_0^2} \right) \frac{1}{n^2} \quad (1)$$

Therefore the energy one gets for a transition from an initial level n_i to a final level n_f is $\Delta E = E_{n_i} - E_{n_f}$.

This energy is then set free in the form of light and one can relate it to the wavelength via

$$\frac{1}{\lambda} = \frac{f}{c} = \frac{\Delta E}{hc} = \frac{E_{n_i} - E_{n_f}}{hc} = R \left(\frac{1}{n_f^2} - \frac{1}{n_i^2} \right) \quad (2)$$

where R is called the Rydberg constant, which is given by

$$R = \frac{2\pi^2 me^4 k_e^2}{h^3 c} = 1.0974 \times 10^7 \text{ m}^{-1} \quad (3)$$

2.2 Resolving a Light Spectrum with a Grating

A grating is a collection of small parallel slits. In our case the slits are so fine that we cannot see them any more. A grating has the nice property that light of different wavelength gets diffracted in different directions. Therefore one can resolve a continuous light spectrum (containing light of various wavelengths) into its components.

How does a grating do this trick? To understand this it is sufficient if we look at only two slits with a distance d between them. (d is also called the lattice constant)

We look at two light rays coming from the two slits. If the two rays travel perpendicular to the two slits the two waves will always be in phase, since they travel the same distance. But if the two rays propagate at an angle θ to the normal the two rays will have to travel different distances to reach the eye. (we always assume that the two rays are parallel, which is a good approximation since the eye is almost infinitely far away compared to the distance between the two slits.)

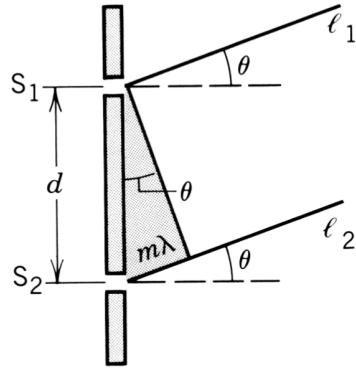


Figure 1: Geometry of the Double Slit

As you may convince yourself by looking at the graph, the difference in path between the two rays is

$$\Delta l = d \sin \theta \quad (4)$$

As in the lab about standing waves we can have different possibilities how these two rays add up. They can add up to form a node or they can add up to form an anti node (or something in between). To form a node every mountain shall be made lower and every valley shall be made higher. Therefore each hill from one wave should meet a valley from the other ray and vice versa. To get an antinode and therefore a

maximum in the total motion of the wave every hill should meet a hill and every valley should meet a valley.

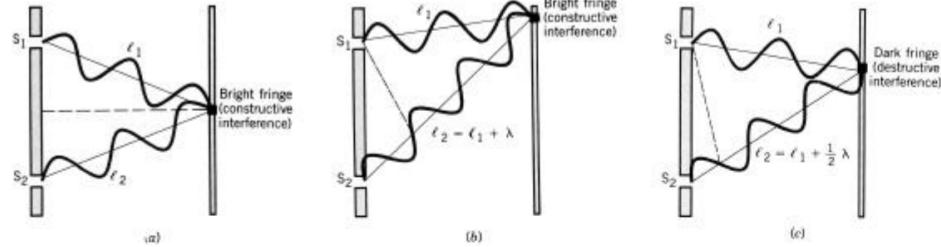


Figure 2: Young's Double Slit

How do we now achieve maximum intensity (node)? We always get this if the two waves are in phase which means that there fits an exact number of wavelength m in the difference between the two waves, i.e.

$$\Delta l = m\lambda \quad (5)$$

m is also called the order. If we now put the two equations for Δl together, we get

$$d \sin \theta = m\lambda \quad (6)$$

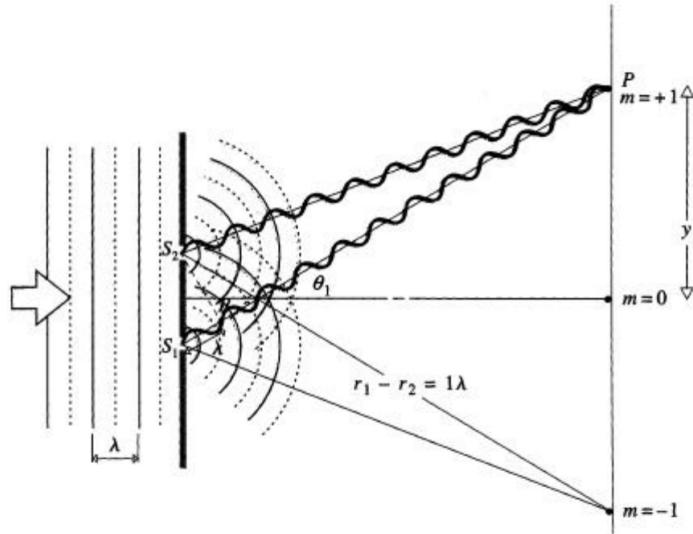


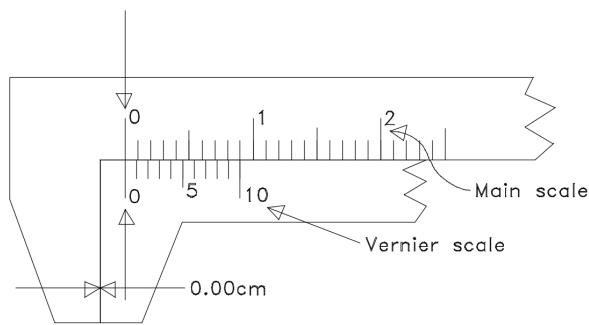
Figure 3: Interference of Light Coming from Two Small Slits

This tells us that given the order m we can determine the wavelength of the light we see through the grating by reading off the angle. This also works the other way

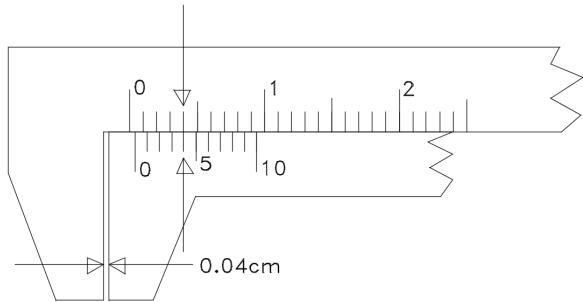
around. (We can also get m by count up the orders starting with 0 at an angle of 0 degrees, and then count how often a particular color repeats as you go to higher or lower angles.)

2.3 How to read a Vernier Scale

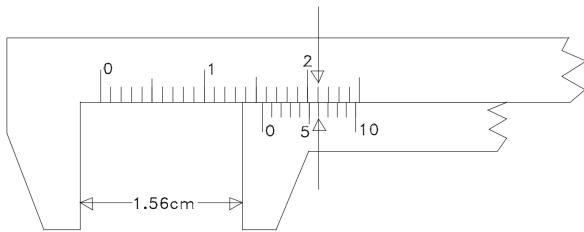
With a usual ruler you can read distances up to a mm resolution. But if you would like to have a measuring device to read with a 0.1 mm resolution you have a problem. In principle you can scratch such a fine scale in the ruler, but in practice you can easily imagine that due to the width of the marks themselves such a scale would be almost impossible to read. That is why we use a trick and introduce a Vernier scale. A Vernier scale works in the following way:



- a) If the zero points of the vernier and main scales are aligned, the first vernier division is $1/10$ of a main scale division short of a mark on the main scale; the 2nd vernier division is $2/10$ of a main scale division short of the next mark on the main scale, etc., and the 10th vernier mark coincides with a main scale mark.



- b) If the vernier scale is now moved to the right until the 4th vernier division is lined up with the nearest main scale division to its right, the distance the 0 point of the vernier scale has moved past the 0 point on the main scale is $4/10$ of a main scale division. Thus the vernier scale gives the fraction of a main scale division that the zero point of the vernier has moved beyond any main scale mark.



- c) In this last example, the zero line of the vernier scale lies between 1.5 cm and 1.6 cm on the main scale. The fraction of a division from 1.5 cm to the zero line can be determined as follows: Line number 6 of the vernier scale coincides with a line on the main scale, so the zero line on the vernier has moved $6/10$ of a main scale division away from the main scale line 1.5. Hence, the reading, to the nearest hundredth, is 1.56. The maximum uncertainty is now ± 0.01 cm, corresponding to the smallest division on the vernier scale. The use of the vernier has reduced the maximum uncertainty of the caliper measurement from ± 0.1 cm (without the vernier) to ± 0.01 cm.

In this experiment we will not use a simple Vernier scale since we do not measure distances, but we will use an angular Vernier scale to measure sub-divisions of angles. If you look at the main scale you will see that the scale is divided in half-degree steps. (You have slightly longer marks for the full degrees and shorter marks for the half degrees.) For angles the next smaller unit is not $1/10$ th of a degree but minutes, which is $1/60$ th of a degree. But since we can already read the scale up to $1/2$ degree we need not have a Vernier scale with 60 divisions but with 30 divisions.

So you first use the 0 mark on the Vernier scale to read off the angle in degrees and if you are in between 0 and 30 minutes or 30 and 60 minutes. Then you look which of the marks on the Vernier scale exactly matches a mark on the other scale. So e.g. if the 0 mark on the Vernier scale is between 20.5 and 21 degrees and the 13th mark on the Vernier scale matches a mark on the degree scale we would have a reading of $20 + 1/2 + 13/60$ degrees or 20 degrees and 43 minutes.

3 Experiments

3.1 Adjusting the Spectrometer

The first part of the experiment will basically include the procedure to set up the equipment. This should be done with as much care as possible. Only then we will be able to measure the wavelength on the limit of our apparatus. If you don't set up the spectrometer correctly you will get systematical errors, which screw up your data.

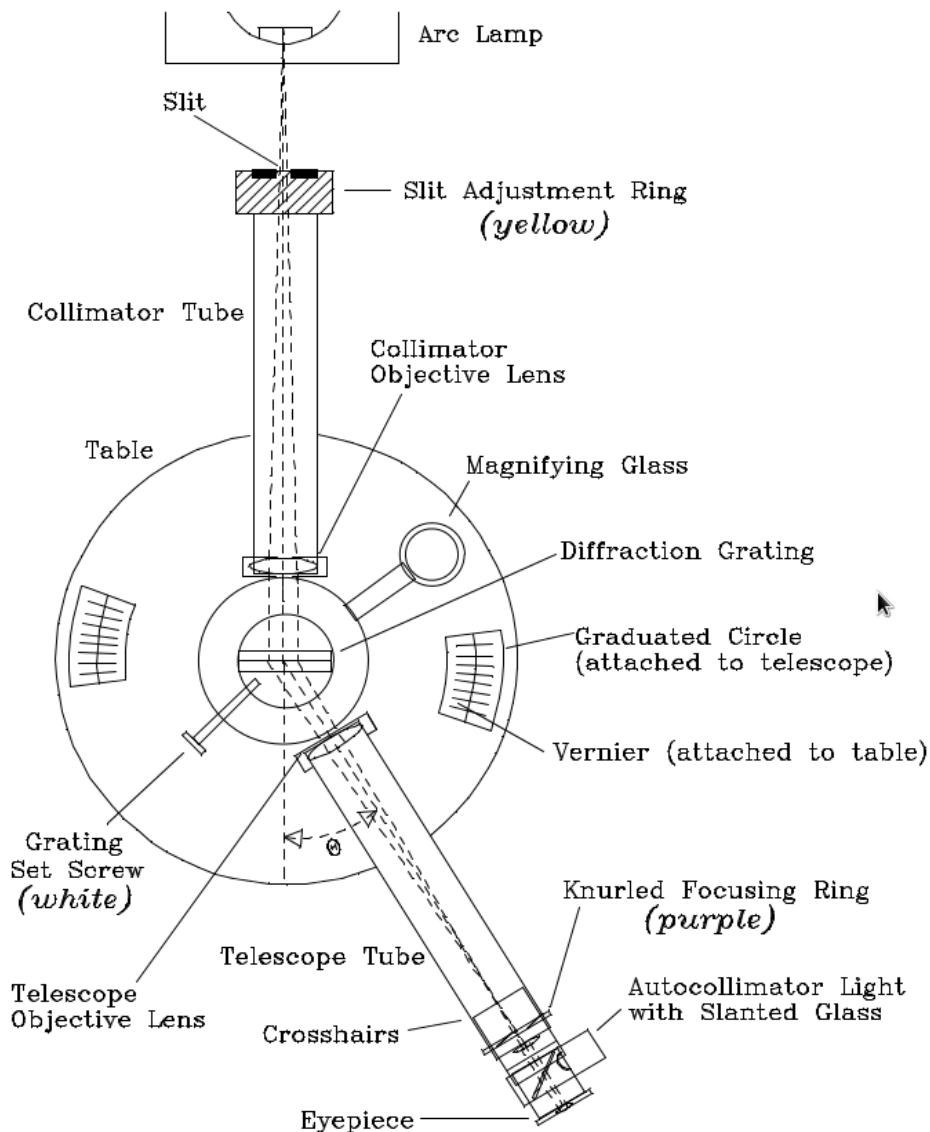


Figure 4: Schematic of the Spectrometer

3.2 Obtaining the Lattice Constant

After adjusting the spectrometer we will measure the yellow line of a Helium discharge lamp. Since we know that the wavelength of this light is $\lambda = 5.8756 \times 10^{-7}$ m, we can determine the lattice constant of the grating quite accurately. Even though the lattice has written on it 600 lines/mm (which is only an approximate value anyway), we want to get the lattice constant with 5 relevant digits and not just 3, and this means that we have to measure it!

3.3 Measuring the Spectrum of Hydrogen Atoms

With the lattice constant determined in the previous part we now measure the wavelength of the light emitted from the Hydrogen discharge lamp.

4 Step by Step List

4.1 Adjusting the Spectrometer

- Take the grating out off the holder and close the green knob.
- Rotate the yellow knob such that the slit is about half open.
- Look through the eyepiece and turn the purple focusing ring until you see a sharp image of the slit.
- Loosen the red knob and move the telescope tube until the cross hair is in the middle of the slit. Tighten red knob.
- Now open the green knob and turn the tabletop such that the 0 mark from the Vernier scale with the magnifying glass is lined with either the 180 or the 360 from the outer scale. (Always use only this Vernier scale and don't switch to the other one in between)
- Close the green knob (and don't open it again for the rest of the experiment!).
- Now you can fine adjust the relative position of the inner and outer scale by turning the blue knob. Make sure that the lineup between the 0 on the Vernier and the 180/360 is done as careful as possible. (For some of the spectrometers there is a small mark to the left of the 0 mark on the Vernier scale. Make sure that you line up the 0 mark and not this extra mark with the 180/360.)
- Now put the grating in the holder such that it is perpendicular to the telescope tube-collimator tube line. Close the white screw to lock the grating.

4.2 Obtaining the Lattice Constant

- Switch on the Helium lamp and line the spectrometer up such that you can see the slit well illuminated by the lamp as you look through the spectrometer.
- Put the black cardboard over the front end of your collimator tube and you can use the black piece of cloth over the spectrometer to block light from the surrounding. Also this part and most important the next one should be performed

in the dark, therefore switch off the light and use the lamps provided (they are supposed to be that dim) to read the scale.

- Open the red knob and move the telescope tube to the left until the crosshair is in the center of the yellow line. (You should first see a few blue and green lines, then the isolated yellow line and then red lines. The yellow line should be somewhere around 20 degrees)
- Note down the angle in degrees and minutes at which you see the 1st order of the yellow line. Do the same on the right side and average these two numbers.
- Use the average and plug it into the grating equation ($m = 1$) to determine the lattice constant d (at least 5 relevant figures).
- How many lines per mm does this lattice constant correspond to?
- Can you also see the 2nd order and 3rd order yellow lines on either side?

4.3 Measuring the Spectrum of Hydrogen Atoms

- Now switch off the helium lamp and set up the spectrometer for the hydrogen lamp.
- The light purple line you see in the middle is the 0th order and not yet one of the lines we are going to measure.
- There are 4 visible lines for the Hydrogen spectrum. One red (furthest out), one greenish blue, one purple blue and one dark purple. The dark purple line is very faint and you may not be able to see it, so don't despair.
- Measure the angles of the 4 lines on both sides and average the angles.
- How many orders do you see on either side? (look e.g. for the red line and count how often it appears as you go further out)
- Do the higher orders overlap? (i.e. does a new order start before the old one ended?)

4.4 Analyzing the Data

- Put together formulas (2) and (6) and solve the resulting formula for n_i .
- Use your data and the value for R from above and $n_f = 2$ to determine n_i .
- n_i should be an integer number labeling from which initial atomic shell the electrons in the hydrogen atoms fell to the 2nd atomic shell.

- Are the results integers (or close to them)?
- Are the numbers (integers) the one they are supposed to be?
- For which color did the electrons jump from the lowest shell? Explain why you could have predicted that anyway!
- Within how many % have you been able to measure the data? Compare that to other experiments you have performed in this course!
- Discuss briefly the factors that determine the uncertainty in your measurements. Which of them are random, which systematic?

5 Lab Preparation Examples

Hydrogen Spectrum:

1. What is the wavelength of emitted if an electron in an hydrogen atom makes a transition from $n_i = 3$ to $n_f = 1$?
2. For a hydrogen atom give all the transitions that fall within the visible range of the light spectrum (i.e. between $\lambda_{\min} = 400 \text{ nm}$ and $\lambda_{\max} = 800 \text{ nm}$). Give the transitions in the following form: A transition from energy level $n_i = 3$ to $n_f = 1 : 3 \rightarrow 1$.

Hint: It may be simplest to calculate the energies $\lambda_{\min} = 400 \text{ nm}$ and $\lambda_{\max} = 800 \text{ nm}$ correspond to. Then make a list of E_n and choose the differences between the energy level that fit in the calculated range.

3. What wavelengths do the transitions in the previous problem correspond to?
4. If you look at the light of an object at high temperature (e.g. a normal light bulb), you will find that it usually emits a continuous spectrum without any gaps or lines in it. (This is due to the effect that the atoms in a solid (or plasma) are so close together that they disturb each other so strongly that the initially sharp transition lines of the atoms get so broadened that they appear as a continuous spectrum.)

But if you now look at the spectrum of the sun (or another star), you will find a continuous spectrum that has a number of sharp black lines, so no (or almost no) light of that frequency reaches you; the light of that frequency is missing in the spectrum. These black lines are at exactly the same positions where you would see the colored lines in emission spectra (that is what we do in lab) of the elements present on the sun. For example since hydrogen is present on the sun we will find black lines at the wavelength calculated in 3.

Can you give a short explanation of how we get these inverse spectra (compared to the ones we see in the lab from emission spectra) from the sun?

Spectrum with a Grating:

5. What is the lattice constant d if the lattice has 400 lines per mm?
6. If you have a grating with 600 lines per mm at what angle would you observe the $m = 1$ maximum for a wavelength of 600 nm?
7. If you have $d = 10^{-1} \text{ m}$ and $\lambda = 1 \text{ cm}$. at what angles will you see the maxima for $m = 1, 2, 3$?

Experiment 2-10

Absorption of Beta and Gamma Rays

Objective: To study the behavior of gamma and beta rays passing through matter; to measure the range of beta-particles from a given source and hence to determine the endpoint energy of decay; to determine the absorption coefficient in lead of the gamma radiation from a given source.

Safety Note: The radioactive sources used in this experiment are extremely weak and therefore can be handled without any danger of overexposure to radiation. However, it is always prudent not to handle any radioactive source more than is necessary.

1 Discussion

All nuclei heavier than lead (and many isotopes of lighter nuclei) have a finite probability of decaying spontaneously into another nucleus plus one or more lighter particles. One of the decay products may be an alpha-particle (two protons and two neutrons—the stable nucleus of a helium atom). Alternatively, a nucleus with more neutrons than it can maintain in stability may decay by emission of an electron from the nucleus (beta-decay) which corresponds to the conversion of a neutron to a proton. These electrons may emerge with a kinetic energy of up to several MeV.

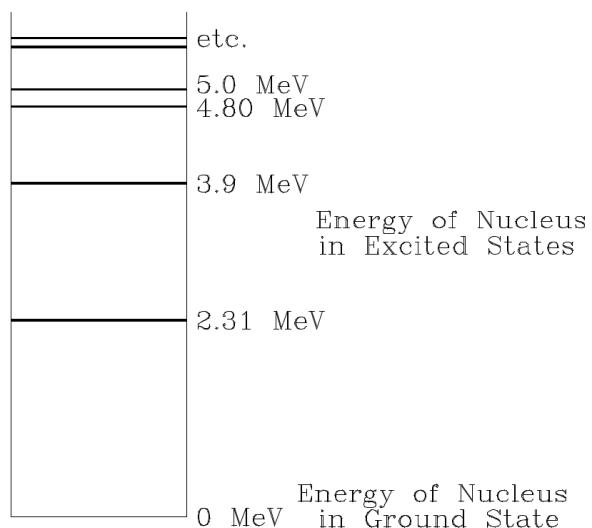


Figure 1: Typical Nuclear Energy Level Diagram.

After alpha or beta decay, the residual nucleus may be left in an excited state. In this case, a transition to a state of lower energy of the same nucleus will occur almost immediately with the emission of a photon (gamma ray). The spectrum of photons emitted from the various excited states of a nucleus will have discrete frequencies ν , corresponding to transitions $\Delta E = h\nu$, between discrete energy levels of the nucleus. The gamma ray spectrum from an excited nucleus is thus **analogous** to the line spectrum of visible radiation from an atom due to excited electrons, with the notable difference that the MeV energy changes of the nucleus are approximately 10^6 times as large as energy changes in transitions between atomic states (where $\Delta E_{\text{atomic}} \approx$ several eV). See Figure 1.

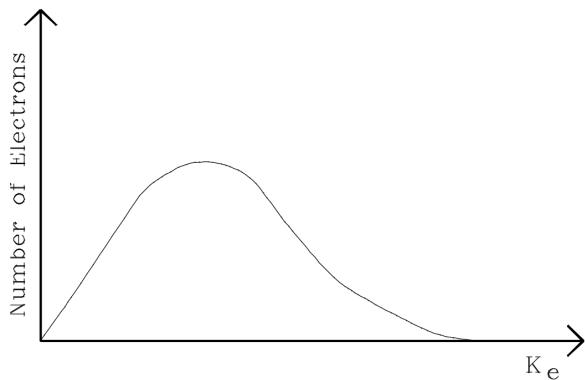


Figure 2: Beta Ray (Electron) Energy Spectrum

In early experiments on beta-decay, it was observed that each decay was not a simple one in which an electron and the recoil nucleus came off with equal and opposite momentum. The electrons, in fact, were emitted with a continuous spectrum of energies (Figure 2). It was subsequently suggested by Pauli and Fermi that in each decay, another particle of zero mass and charge, called the neutrino, was emitted. Experimental verification of the neutrino has been obtained by observation of its rare interaction with matter.

2 Detection of Charged Particles (the Geiger Counter)

When an energetic charged particle traverses matter, it will give electrostatic impulses to nearby electrons and thus ionize atoms of the material. Most methods for detecting nuclear particles rely on observing the results of this ionization. In a photographic emulsion, a trail of ionized grains will show up when the film is developed. A trail of droplets or bubbles forms along the wake of ionization left by a charged particle passing through a cloud or bubble chamber. In a scintillation counter, ionized molecules will very rapidly radiate visible light. A charged electroscope will be discharged by ion

pairs created in air. In the **geiger counter**, which will be used as a detector in this experiment, the ionization produced by a charged particle causes a violent electrical discharge.

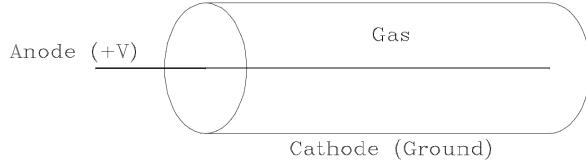


Figure 3: The Geiger Counter.

As shown in Figure 3, the counter consists of a metal cylinder (cathode) with insulating ends supporting a fine axial wire (anode). When a charged particle ionizes the gas (e.g., argon) in the tube, electrons will move toward the positively charged anode wire, accelerated more and more by the rapid increase in electric field near the wire. When an electron acquires kinetic energy greater than the ionization energy of the gas molecules, it can create by collision a new ion and electron, which in turn can accelerate and create another ion, etc., thus initiating an avalanche of charge. The process, called the **Townsend avalanche**, is possible only if the voltage maintained between anode and cathode is sufficiently high.

The basic counter circuit, shown in Figure 4, supplies a positive high voltage of up to 900 volts to the center wire. When an avalanche occurs, current flows through R , the counter side of R drops in potential, and this negative pulse is fed through C to a stage of amplification and then to a scaling device.

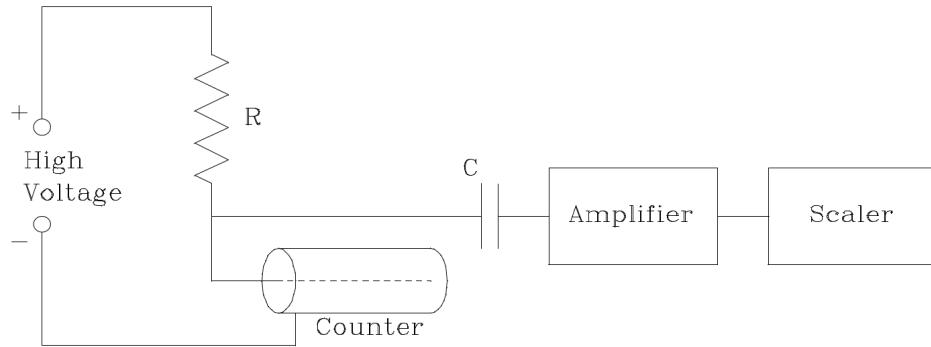


Figure 4: Geiger Counter Circuit.

For a large enough value of R , the voltage drop across R would be sufficient to stop the discharge after a count has been recorded. This method of quenching the discharge, however, has the disadvantage of creating a long time constant, $\tau = RC$, for reestablishing voltage across the counter, i.e., of creating a long “dead time” during which the passage of another particle cannot be detected. The quenching can be

achieved for a lower value of R by adding halogen vapor to the gas to help absorb ion pairs. Nevertheless, the resultant dead time is substantial and must be taken into account when the particle flux is high.

3 Energy Loss and Range of Beta Particles

Because of its ionizing action (Figure 5), a *charged*, incident particle in matter will continuously lose kinetic energy, and the particle will subsequently come to rest after traversing a path length called its **range**. *For a particle of known charge and mass, there will be a unique range associated with each incident energy.* A formula can be theoretically deduced for the rate of energy loss—and hence the range—of a particle (of known mass, charge, and initial velocity) in a particular “stopping” material (of known electron density and ionization potential).

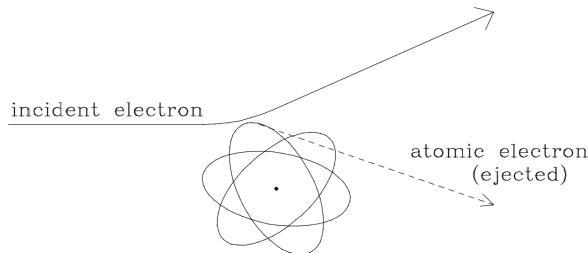


Figure 5: Ionizing Action of Incident Electron.

In each interaction with atomic electrons, however, an incident electron may be scattered through small or relatively large angles, and as it traverses the material it may follow a rather tortuous, winding path (especially at low energies). Therefore, the actual path of the electron may be considerably longer than the observed distance that it penetrates into the material. For this reason, the incident electron range is not sharply determined, and the theoretical calculation is of limited usefulness for electrons of less than 1 MeV energy.

In this experiment, therefore, we will use an approximate empirical relationship between range and energy for low energy electrons:

$$r = \frac{0.412 \text{ g/cm}^2}{\rho} E^{1.29} \quad (1)$$

where r is in cm, E is in MeV, and ρ is the density of the stopping material in g/cm^3 . Note that the density of Aluminum is 2.702 g/cm^3 .

(This result is described in: L. Katz and A. S. Penfold, “Range-Energy Relations for Electrons and the Determination of Beta-Ray End-Point Energies by Absorption,” Revs.Modern Phys. **24**, 1 (1952).)

4 Absorption of Gamma Rays

Gamma rays, or high-energy photons, can interact with matter by three distinct processes:

- 1) Compton Scattering: This refers to a photon-electron collision in which the energy lost by the scattered photon is given to the recoil electron.
- 2) Photoelectric Effect: The photon is absorbed by the atom as a whole, releasing an electron with kinetic energy equal to $E_\gamma - E_b$, where E_γ is the photon energy and E_b is the relatively small binding energy of the electron in the shell from which it is released.
- 3) Pair Production: If the photon has energy greater than 1.02 MeV, it can create an electron-positron pair in the neighborhood of a nucleus. The radiative source used in this experiment does **not** emit photons with energy greater than 1 MeV.

The probability of each of the three processes taking place in a given thickness of material depends on the energy of the photon and the atomic structure of the material. The *total* probability for interaction of photons in Pb, i.e., the sum of the probabilities of the three processes, varies with photon energy as indicated in Figure 6. The ordinate plotted on the graph is μ , the **total linear absorption coefficient** in units of cm^{-1} . It is defined by the equation:

$$\frac{dN}{dx} = \mu N \quad (2)$$

where N is the number of incident photons and dN is the number removed from the beam (i.e. absorbed) in an absorber of thickness dx (in cm). Note that dN and dx are the calculus equivalents of infinitesimally small values of ΔN and Δx , respectively. As in any process where the rate of decrease is proportional to the number present (such as the discharge of a capacitor), the solution of this differential equation is:

$$N(x) = N_0 e^{-\mu x} \quad (3)$$

where $N(x)$ is the number of photons passing through x cm of absorber and $N_0 = N(x)$ at $x = 0$, and e is the base of natural logarithms.

5 Procedure

5.1 Adjustments and Measurement of Errors in Counting

5.1.1 High Voltage Variations

Every Geiger tube that is in good working order has a plateau region in which its counting rate is relatively insensitive to changes in the high voltage supply. This region

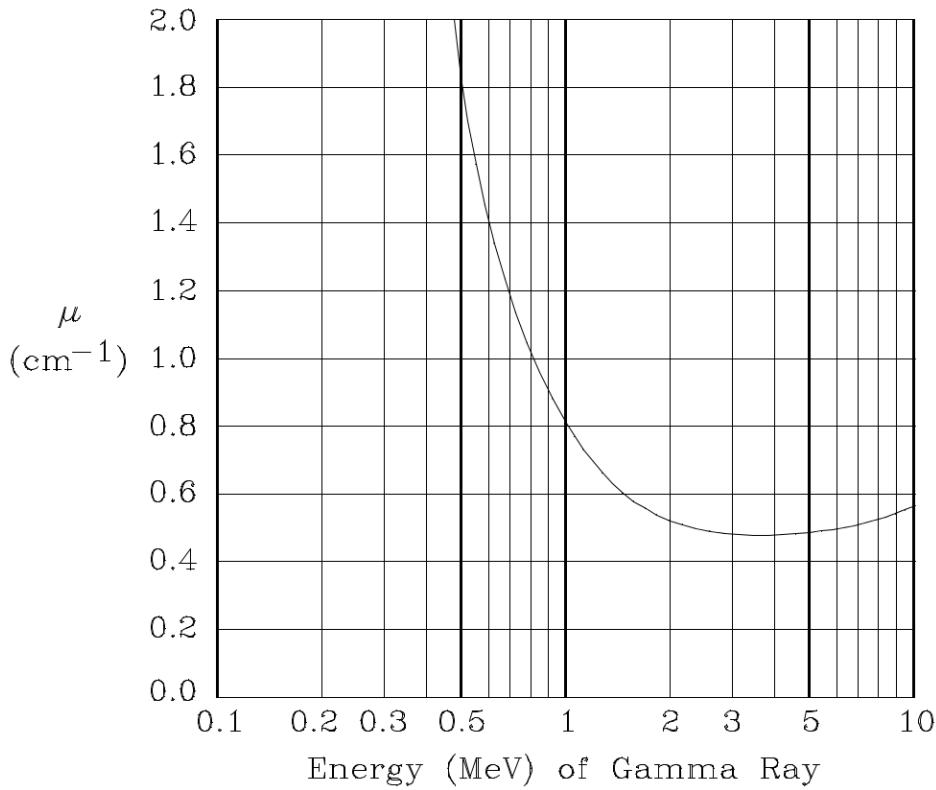


Figure 6: Linear Absorption Coefficient for gamma rays in lead as a function of energy.

can be found by putting a source under the tube and increasing the high voltage from its lowest value until the tube just begins to count. One can then take 15 second counts while raising the voltage in 50 volt steps. The curve of counting rate vs. voltage levels off, i.e., the count rate rises by less than 10% for a 100-volt increase. Do not raise the voltage further as this may damage the tube due to continuous breakdown. Set the high voltage to a value on this plateau for the remainder of the experiment. If this procedure is followed correctly, high voltage variations may be ignored as a source of error.

5.1.2 Statistical Accuracy

Particles decay randomly in time from a radioactive source (over a period short compared to the half-life). The probability distribution for measuring a given number of these counts in a given time interval is an almost bell shaped curve (a Poisson Distribution) centered around N_0 , the most probable value. The distribution will have a standard deviation about the peak of $\pm\sqrt{N_0}$. If N counts are measured in an interval, the best estimate of the error is $\pm\sqrt{N}$.

Note that the magnitude of the statistical error—your uncertainty in the measurement—

increases significantly for trials involving a very small number of counts. For a high counting rate in a given interval, say, 900 counts in one minute, the estimated error will be 30 counts per minute or 3.3% error. However, for a much lower counting rate in the same interval, say 25 counts, the error of ± 5 counts per minute amounts to a 20% statistical error. In order to achieve the same precision as in the first case, it would be necessary to collect 900 counts—in other words, to take a 36-minute measurement. While such a long measurement is impractical for this lab, you should keep in mind the relationship between a small number of counts and higher errors, and do your best to minimize these errors by taking longer measurements when necessary.

5.1.3 Measuring Background Radiation

In order to make accurate counting measurements of the sources, it is necessary to know the counting rate due to natural background radiation (mostly cosmic rays coming through the earth's atmosphere). Additionally, there will be some excess counts due to the Cs-137 gamma sources nearby in the room, whose gamma rays can pass through the side of the detector. At a distance of 30 cm, for example, the Cs-137 source contributes roughly as many counts as the natural background radiation (doubling the distance would reduce its contribution to one-fourth this level). It's best to try to minimize these secondary effects—by keeping your detector far from others' sources, and by shielding your own Cesium source with the lead sheets when not in use), but it's even more important to try to keep these background effects *constant*. If all of your data is shifted by roughly the same constant amount, then it is possible to isolate the results you're interested in by subtracting out this constant background.

5.2 Range of Beta Particles

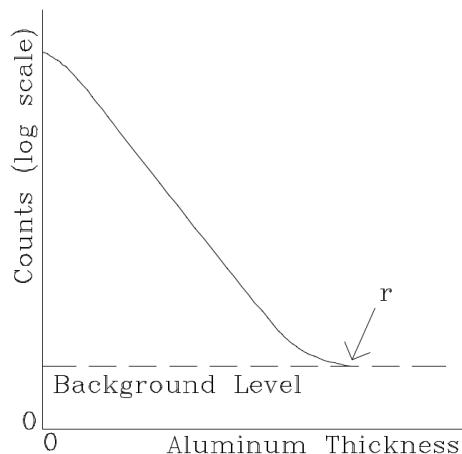


Figure 7: Beta Absorption Curve

Thallium 204 beta decays to Lead 204 with a half-life of 3.9 years. The range of the most energetic of the decay electrons can be determined by placing aluminum foil absorbers between the source and the geiger counter.

A typical absorption curve is shown in Figure 7. The maximum range r is the point where the absorption curve meets the background. You should start by making a careful measurement of background, and you should repeat this measurement after taking the absorption curve to check for constancy.

Place the Thallium source on the second shelf below the detector, as shown in Figure 8, to maximize the number of counts while leaving enough room to stack aluminum absorbers. Begin taking measurements for the absorption curve, adding aluminum foils until the counting rate reaches the background level. (*Note that the thickness of the aluminum absorbers is marked on the foil in **mils**, or thousandths of an inch, not millimeters.*)

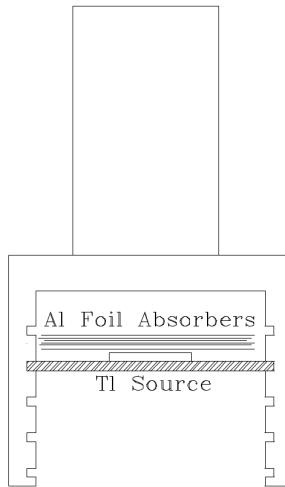


Figure 8: Beta Absorption

Make a table of the results, including the background level, and include estimates of the statistical error in each measurement. Note that you *should not* subtract background from your data for this experiment. Plot the results on the semi-log paper provided (if your measurements are of different durations, you will need to scale them to a constant duration which is convenient for graphing).

Determine the approximate value of the maximum range from the graph, and use Eq. (1) to compare your result with the value of 0.765 MeV for the maximum beta energy for Thallium 204 as measured in a magnetic spectrometer.

5.3 Absorption of Gamma Rays

The source for this part of the experiment is Cesium 137, which decays with a half-life of 30 years to Barium 137 with the emission of a 0.52 MeV beta ray (Figure 9). The resulting Barium is in an excited state and decays by emitting a 0.662 MeV gamma ray almost instantaneously. Lead absorbers are used for the gamma absorption study. They are thick enough(0.062 inches) so that one absorber will stop all the betas.

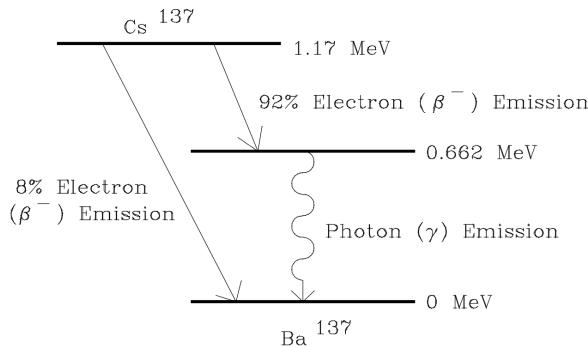


Figure 9: Cesium 137 Decay Scheme

The gamma rays are detected by means of the same geiger counter used in the previous part. Note that the efficiency of the geiger counter for detecting photons is much less than for detecting the beta rays, since it depends on a collision of the photon with the gas or wall of the counter, resulting in the emission of an electron, which in turn initiates the discharge.

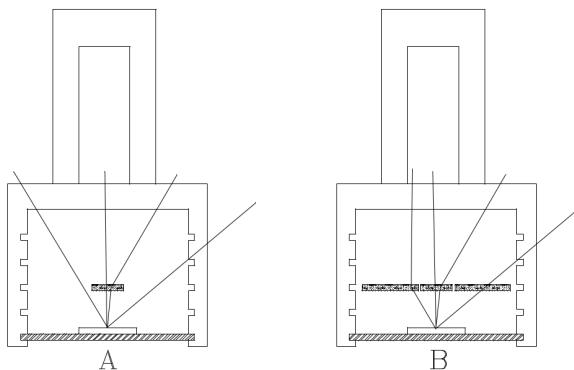


Figure 10: Scattering Effects: In (A), the gamma ray on the left passes outside of the detector tube. In (B), it can be seen how increasing the area of lead absorber can cause the same gamma ray to scatter *into* the detector.

Figure 10 represents one of the effects of gamma ray scattering in the lead sheets. Increasing the *area* of lead through which the gamma rays pass tends to *increase*, rather

than decrease, the number of counts one measures, since gamma rays which otherwise would not have entered the detector may now be scattered into it.

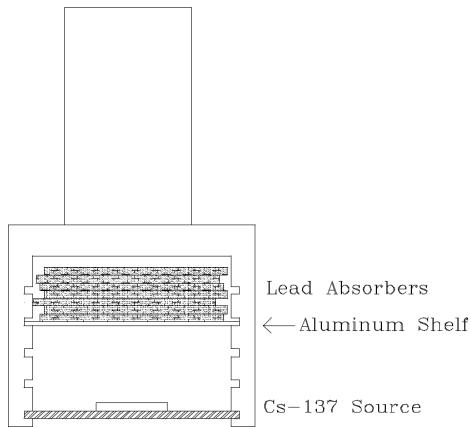


Figure 11: Absorption of Gamma Rays Set-Up.

The arrangement shown in Figure 11 is designed to reduce the effect of such scattering. By keeping the lead sheets high above the source, one reduces the excess area exposed to gamma rays, and one reduces the effective difference in area between the top and bottom sheets as well.

Take measurements for an absorption curve. Note how this differs fundamentally from beta absorption: there is no maximum range for gamma rays passing through lead; rather, one expects to lose a fixed *fraction* of the remaining gamma rays passing through each successive layer of lead absorber. For this reason, the theoretical absorption curve never intersects the background curve. When the absorption curve is plotted on semi-log graph paper with *background subtracted*, the exponential decay curve appears as a straight line.

Make a table of the data and estimated errors. Plot the data on semi-log graph paper as in Part II.

Using a straight line fitted to the data, find the absorption coefficient μ by choosing two points on the line a distance x apart, substituting their y values for N_0 and $N(x)$ in Eq. (3) and solving for μ . Use this value and Figure 6 to estimate the energy of the gamma rays, and compare with the accepted value. What factors limit the precision of this estimate?