# Mirela Nicolov Zoltan Szabadai 

## Watz Claudia

## PHARMACEUTICAL PHYSICS



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## Editura "Victor Babeş"

Piața "Eftimie Murgu" nr. 2, cam. 316, 300041 Timișoara
Tel./Fax 0256495210
email: evb@umft.ro
www.umft.ro/editura

# Director General: Prof. univ. emerit Dr. Dan V. Poenaru 

Colecția: Ghiduri şi îndrumătoare de laborator
Coordonator colectie: Prof. univ. Dr. Adrian Vlad

Referent științific: Prof. univ. dr. Codruța Șoica
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## LABOR PROTECTION RULES IN THE PHARMACEUTICAL PHYSICS LABORATORY

## General rules

- During the hours of practical work, a protective gown will be worn and synthetic clothing will be avoided
- Only the presence of the experimental instruments, the workbook and the computer is allowed on the work table. Personal items will be stored in places designated for their storage
- Practical works will be carried out only under the supervision of the teacher
- Before using substances check the label (three times)
- The ongoing experiment will not be left unattended
- It is strictly forbidden to taste the chemicals from the laboratory
- The use of laboratory glassware for the purpose of drinking or eating is strictly prohibited
- Glassware will be checked to be in good working condition and washed before and after use (three times with distilled water)
- In order to avoid contamination of substances, it is not allowed to exchange stoppers between containers, nor to insert pipettes from one container to another, without first washing it


## In case of chemical accident

In case of a burn with corrosive substances, it is necessary to wash the wound with plenty of water and treat with various substances, depending on the chemical agent that caused the burn:

- burns caused by acids : a base will be used for neutralization - $2 \%$ sodium bicarbonate solution -burn caused by bases: an acid will be used for neutralization - $2 \%$ boric acid solution
- thermal burn: apply special pharmaceutical preparations addressed for the treatment of burns eg: OXIMED A medical visit is mandatory after any accident, for all the people who were involved in it !


## PHYSICAL QUANTITIES AND UNITS OF MEASUREMENT

Physical quantities define properties of their bodies characterize processes in which the changes that occur can be described quantitatively.

Examples of physical quantities are: mass, temperature, speed, electric charge.
Physics has long been called the science of measurement, because the study of physical phenomena involves the measurement of the quantities that characterize them.

Measurement is a process by which it can be compared the respective physical quantity with a welldefined quantity, of the same nature, which a was chosen as the unit of measure. This comparison (or measurement) is done by means of a measuring instrument.

## EXAMPLES OF UNITS OF MEASUREMENT:

1 meter for length, 1 second for duration, 1 kg for mass .
Some physical quantities are fundamental quantities, being defined only by
description of the measurement procedure.
For example, distance is determined by measurement with a ruler, and time by measuring with a clock. Other physical quantities are derived quantities, as they are obtained through calculation formulas that use the fundamental quantities.

For example, speed represents the ratio between the distance traveled by a body in a period of time. Over time, different systems of measurement units have been used such us, fundamental/basic quantities and also other corresponding units of measurement.

Today, the International System of Units, SI, is used most frequently, employing the following fundamental sizes and units of mesurement:

| No. crt | Physical quantity | Unit of <br> measurement | SI | CGS |
| :---: | :--- | :--- | :--- | :--- |
| 1 | Time | Second (s) | $[\mathrm{t}]_{\mathrm{SI}}=1 \mathrm{~s}$ | $[\mathrm{t}]_{\mathrm{CGS}}=1 \mathrm{~s}$ |
| 2 | Distance | Meter (m) | $[\mathrm{d}]_{\mathrm{SI}}=1 \mathrm{~m}$ | $[\mathrm{~d}]_{\mathrm{CGS}}=1 \mathrm{~cm}$ |
| 3 | Mass | Kilogram (kg) | $[\mathrm{m}]_{\mathrm{SI}}=1 \mathrm{~kg}$ | $[\mathrm{~m}]_{\mathrm{CGS}}=1 \mathrm{~g}$ |
| 4 | Temperature | Kelvin | $[\mathrm{T}]_{\mathrm{SI}}=1 \mathrm{~K}$ | $[\mathrm{~T}]_{\mathrm{CGS}}=1 \mathrm{~K}$ |
| 5 | Amount | Mole | $[\mathrm{v}]_{\mathrm{SI}}=1 \mathrm{~mol}$ | $[\mathrm{v}]_{\mathrm{CGS}}=1 \mathrm{~mol}$ |
| 6 | The intensity of the | Amps | $[\mathrm{I}]_{\mathrm{SI}}=1 \mathrm{~A}$ | $[\mathrm{I}]_{\mathrm{CGS}}=1 \mathrm{~A}$ |
| 7 | Luminous intensity | The candle | $[\mathrm{E}]_{\mathrm{SI}}=1 \mathrm{~cd}$ | $[\mathrm{E}]_{\mathrm{CGS}}=1 \mathrm{~cd}$ |

METER - UNIT OF LENGTH (m) The meter is the length equal to 1650763.73 vacuum wavelengths of the radiation that corresponds to the transition between the 2 p10 and 5 d 5 energy levels of the Krypton 86 atom.

## Observation

The definition of the meter was adopted by the 11th General Conference of Weights and Measures in 1960 through the 6th resolution.

THE KILOGRAM - UNIT OF MASS (Kg): The kilogram is the mass of the international prototype. The term weight denotes the same quantity with a force $G=\mathbf{m} \cdot \mathbf{g}$ - where $\mathbf{g}$ is the gravitational acceleration; $g=9.80665 \mathrm{~m} / \mathrm{s}^{\mathbf{2}}$

## Observation:

The definition of the kilogram was adopted by the first and third General Conferences of Weights and Measures in 1889 and 1901, respectively. The kilogram as an "international prototype" is kept at the International Bureau of Weights and Measures in Sevres, under the conditions set by the General Conference of Weights and Measures from 1889.

SECOND - TIME UNIT (s) The second is the duration of 9192631770 periods of the radiation that corresponds to the transition between the two hyperfine levels of the fundamental state of the Cesium 133 atom.

Note:
The definition of the second was adopted by the 13th Conference of Weights and Measures from 1967 through the 1st resolution.

AMPERE - THE UNIT OF INTENSITY OF ELECTRIC CURRENT (A) The ampere is the intensity of an electric current which, maintained in two straight parallel conductors of infinite length and with negligible circular section, placed in a vacuum, at a distance of 1 m from each other, will produce between these conducting a force of 1 N over a length of 1 m .

Observation:
The definition of the Ampere was adopted by the 9th conference of Weights and Measures from 1948 through the 2nd resolution.

CANDELA - THE UNIT FOR LUMINOUS INTENSITY (Cd) The candela is the luminous intensity in a given direction of a source that emits monochromatic radiation with a frequency of $540.10{ }^{12} \mathrm{~Hz}$ and whose energy intensity in this direction is

## Observation:

The definition for the candle was adopted by the 16th Conference of Weights and Measures from 1979

THE KELVIN - UNIT OF THERMODYNAMIC TEMPERATURE (K) A Kelvin is the fraction $1 / 273.16$ of the thermodynamic temperature of the triple point of water.

Observation:

The Kelvin definition was adopted by the 13th Conference of Weights and Measures in 1967 through the 4th resolution.

The same unit of measurement and the same symbol are used to evaluate a temperature range.

## THE MOLE-UNIT FOR THE QUANTITY OF SUBSTANCE (mol)

a) The mole represents the amount of atoms determined experimentally to be found in 0.012 Kg of C12.
b) Whenever the mole is used, the elementary entities must be specified, they can be atoms, molecules, ions, electrons, other particles or groups specified by such particles.

## Observation:

The definition of the Mole was adopted by the 14th Conference of Weights and Measures from 1971 Other units of measure:

The plane angle is measured in radians .
The radian is the plane angle with the vertex at the center of a circle that delimits on the circumference of the circle an arc whose length is equal to the radius of the circle.

The definition of the radian was adopted by the International Organization for Standardization (ISO) through recommendation R31, part 1.

The solid angle - is measured in steradians .
The steradian is the solid angle with the apex in the center of a sphere that delimits on the surface of the sphere an area equal to the area of a square whose side is equal to the radius of the sphere.

The definition of the radian was adopted by the International Organization for Standardization (ISO) through recommendation R31, part 1.

The first length standard was a bar made of a platinum-iridium alloy called a standard meter, in 1959 the Yard is defined: 1 yard $=0.9144 \mathrm{~m}$;

1 inch $=2.54 \mathrm{~cm}$

## OTHER UNITS OF MEASUREMENT:

## Units of measurement for force :

1 dyn $=10^{-5} \mathrm{~N}$
$1 \mathrm{kgf}=9.81 \mathrm{~N}$

## Units of measurement for mass:

1 pound $(\mathrm{lb})=453.5923 \mathrm{~g}=0.4535923 \mathrm{Kg}$
1 ounce $(\mathrm{oz})=28.35 \mathrm{~g}$
Units of measurement for energy
$1 \mathrm{erg}=10^{-7} \mathrm{~J}$
$1 \mathrm{kcal}=4,186.10^{3} \mathrm{~J}$
$1 \mathrm{eV}=1,602.10^{-19} \mathrm{~J}$
$1 \mathrm{kWh}=3 \cdot 6 \cdot 10^{6} \mathrm{~J}$

## Units of measurement for length

1 inch $=25.4 \mathrm{~mm}$
$1 \mathrm{yd}(\mathrm{yard})=0.9144 \mathrm{~m}$
1 mile $=1609.344 \mathrm{~m}$
1 nautical mile $=1853.18 \mathrm{~m}$

## Units of measurement for area

1 ar $=100 \mathrm{~m}^{2}$
$1 \mathrm{ha}=10000 \mathrm{~m}^{2}$
1 acre $=4046.86 \mathrm{~m}^{2}$
$1 \mathrm{in}^{2}(\mathrm{sq}$ inch $)=6.452 \mathrm{~mm}^{2}$
$1 \mathrm{yd}^{2}(\mathrm{sq}$ yard $)=0.9144 \mathrm{~m}^{2}$
Units of measurement for volume
1 US gallon $=3.785 \mathrm{l}=3785 \mathrm{~m}^{3}$
1 gallon England $=4.546 \mathrm{l}=4546 \mathrm{~m}^{3}$
1 US barrel for oil = 158.981
1 ounce (liquid) $=28.413 \mathrm{~cm}^{3}=1.7339 \mathrm{in}^{3}$
Units of measurement for temperature
Temperature in degrees Fahrenheit: $\mathrm{tF}=(9 / 5) . \mathrm{tC}+32$
Temperature in Reaumur degrees: $\mathrm{tR}=\mathrm{TK}$-218.52
1 gram (SI) $=0.03527$ ounces
1 feddan $($ Egypt $)=1.038$ acres $=0.42$ ha
$1 \operatorname{ardeb}($ Egypt $)=191.6121$

1 hamlah $($ Egypt $)=74,890 \mathrm{~kg}$
1 drachma $($ Turkey $)=3.207 \mathrm{~g}$
1 cuffisco $($ Sicily $)=22.19811$
1 cuartilo (Mexico)=0.547 1
1 cuarteron (Spain) for oil=0.1511
1 dinero $($ Spain $)=1.199 \mathrm{~g}$
1 chittak $($ India $)=58.319 \mathrm{~g}$
1 buncal (Indonesia) $=46.344 \mathrm{~g}$
1 bercocvic $($ Russia $)=163.85 \mathrm{~kg}$
1 pud (Russia) $=16.381 \mathrm{~kg}$
1 barrel for cement (USA) $=170,550 \mathrm{~kg}$
1 barrel for oil $($ USA $)=158.987 \mathrm{~kg}$
1 iron (mass unit for skin thickness) $=1 / 48 \mathrm{in}=0.539 \mathrm{~mm}$
1 koku $($ Japan $)=180.4857 \mathrm{dm}^{3}$
1 kivan $($ Japan $)=3.75 \mathrm{~kg}$

## MEASURED LENGTHS - EXAMPLES

| Distance | Unit of measurement <br> meters (m) |
| :--- | :--- |
| Distance to the most distant quasar | $6 \cdot 10^{25}$ |
| Distance to the nearest nebula <br> (The Andromeda Sea Nebula) | $2 \cdot 10^{22}$ |
| The radius of our galaxy | $6 \cdot 10^{19}$ |
| Distance to the nearest star <br> (Alpha Centauri) | $4,3 \cdot 10^{16}$ |
| Ray of the Sun | $6,9 \cdot 10^{8}$ |
| Earth Ray | $1,4 \cdot 10^{6}$ |
| The height of a man $10^{0}$ |  |
| The thickness of a book | $4 \cdot 10^{-2}$ |
| The thickness of a page in a book | $1 \cdot 10^{-4}$ |
| The size of a polio virus | $1,2 \cdot 10^{-8}$ |
| The radius of a hydrogen atom | $5 \cdot 10^{-11}$ |
| The effective radius of a proton | $1,2 \cdot 10^{-15}$ |

MEASURED TIME INTERVALS

| Time intervals | Unit of measurement <br> Seconds (s) |
| :--- | :---: |
| Age of the Earth | $1,2 \cdot 10^{17}$ |
| The average lifespan of a human | $2 \cdot 10^{9}$ |
| The duration of the Earth's revolution around the Sun $(1$ <br> year | $3,1 \cdot 10^{7}$ |
| The duration of the Earth's rotation around its axis (1 day) | $8,6 \cdot 10^{4}$ |
| The half-life of free neutrons | $7 \cdot 10^{2}$ |
| The time between two normal heartbeats | $8 \cdot 10^{-1}$ |
| The period of the tuning fork that gives the note A | $2,3 \cdot 10^{-3}$ |
| The typical rotation period of a molecule | $1 \cdot 10^{-12}$ |
| The half-life of neutral pawns | $2,2 \cdot 10^{-16}$ |

For force: $1 \mathrm{dyn}=10^{-5} \mathrm{~N}$ or $1 \mathrm{kgf}=9.81 \mathrm{~N}$
For mechanical work : $1 \mathrm{erg}=10^{-7} \mathrm{~J}$
$1 \mathrm{P}($ poise $)=1$ dyn.s $/ \mathrm{cm}^{2}=0.1 \mathrm{~Pa} . \mathrm{s}$
$1 \mathrm{st}=($ stokes $)=1 \mathrm{~cm}^{2} / \mathrm{s}=10^{-4} \mathrm{~m}^{2} / \mathrm{s}$
1 gauss $=1 \mathrm{Gs}=10^{-4} \mathrm{~T}$
1 oersted $=1 \mathrm{Oe}=(100 / 4 \mathrm{Pi}) .(\mathrm{A} / \mathrm{m})$
1 inch $=25.4 \mathrm{~mm}$
1 yard $=0.914 \mathrm{~m}$
Temperature $\mathrm{T}=5 / 9(\mathrm{tF}-32)+273.15$ where tF is the Fahrenheit temperature.

| Other units of measurement for mass |  |  |
| :--- | :--- | :--- |
| No. | Name | CGS |
| 1 | 1 grain | 0.065 g |
| 2 | 1 scruple (apoth) $(=20 \mathrm{~g})$ | 1.29 g |
| 3 | 1 drachma (apoth) (=60g) | 3.88 g |
| 4 | 1 ounce (avdp) | 1.77 g |
| 5 | 1 ounce (apoth) (480 grains) | 31.10 g |
| 6 | 1 ounce (avdp) (437.5 grains) | 28.35 g |
| 7 | 1 livra or pound (apoth) | 373.24 g |
| 8 | 1 livra or pound (avdp) <br> $16 o z=7000$ grains | 453.59 g |

PREFIXES USED FOR MULTIPLES AND SUBMULTIPLES OF MEASUREMENT UNITS

| $10^{-1}$ | deci | $10^{1}$ | deca |
| :--- | :--- | :--- | :--- |
| $10^{-2}$ | centi | $10^{2}$ | hecto |
| $10^{-3}$ | mili | $10^{3}$ | kilo |
| $10^{-6}$ | micro | $10^{6}$ | mega |
| $10^{-9}$ | nano | $10^{9}$ | giga |
| $10^{-10} \mathrm{~m}$ | pico | $10^{12}$ | tera |
| $10^{-12}$ | femto | $10^{15}$ | peta |
| $10^{-15}$ | yocto | $10^{18}$ | exa |
| $10^{-18}$ | ronto | $10^{21}$ | zetta |
| $10^{-24}$ | quecto | $10^{24}$ | yotta |
| $10^{-27}$ |  | $10^{37}$ | ronna |
| $10^{-30}$ |  |  | quetta |
|  |  |  |  |

## THEORY OF ERRORS

## 1. CALCULATION OF ERRORS FOR DIRECTLY MEASURED QUANTITIES

Measurements can only be made with a limited degree of precision and consequently the real value of a quantity cannot be known.

The results of the measurements represent only an approximate value more or less close to the real one. The impossibility of knowing the real value of the quantities to be measured is explained by committing errors during the measurement process.

ERROR represents the difference between the real value and the measured one. Making errors is inevitable, but they can be reduced by taking certain precautions. It is inevitable to know the degree of precision of the measurements which can be appreciated by evaluating the committed error.

The determination of a physical quantity can be done directly, when the result is read directly from the measuring instrument, or indirectly, when it is obtained by calculating it from other directly determined physical quantities (for example, the length 1 , time t , mass m are read directly, whereas the speed v is obtained from the calculation.)

Errors are classified as systematic errors or accidental errors.
Characteristic of the systematic errors is the distortion of the measurements every time in the same direction, either increasing or decreasing all the measurements.

The accidental errors, which are due to random causes, take different values in both directions, in the successive measurements of the same quantity. They are not equal in absolute value, generally differing in different measurements.

Accidental errors obey the laws of probability calculation and are characterized by the following properties:
$\checkmark$ The probability of committing a positive error is equal to that of committing a negative error.
$\checkmark$ The probability of making a small error is equal to that of making a large error.
$\checkmark$ The sum of accidental errors produced when measuring the same quantity, tends to zero when the number of measurements is very large.

The real error is the difference between the measured value and the real value of a quantity:

$$
\begin{aligned}
& \Delta x_{1}=x_{1}-x \\
& \Delta x_{2}=x_{2}-x \\
& \Delta x_{3}=x_{3}-x \\
& \cdots \cdots \cdots \\
& \Delta x_{n}=x_{n}-x \\
& --------- \\
& \sum_{i=1}^{n} \Delta x_{i}=\sum_{i=1}^{n} x_{i}-n x
\end{aligned}
$$

$$
\text { or } x=\frac{\sum_{i=1}^{n} x_{i}}{n}=\frac{\sum \Delta x_{i}}{n}=\bar{x}=\varepsilon
$$

where $\bar{x}$ - is the average of the values of $\mathrm{x}_{\mathrm{i}}$
is the average of the errors
for a large number of errors: $\varepsilon \rightarrow 0$

## AVERAGE OF ABSOLUTE VALUES OF INDIVIDUAL ERRORS

$$
\begin{aligned}
& \left|\Delta x_{1}\right|=\left|x_{1}-x\right| \\
& \left|\Delta x_{2}\right|=\left|x_{2}-x\right| \\
& \left|\Delta x_{3}\right|=\left|x_{3}-x\right|
\end{aligned}
$$

$$
\left|\Delta x_{n}\right|=\left|x_{n}-x\right|
$$

$$
\text { and } \bar{x}=\frac{\left|\Delta x_{1}\right|+\left|\Delta x_{2}\right|+\left|\Delta x_{3}\right|+\ldots+\left|\Delta x_{n}\right|}{n}
$$

A possibility of choosing individual error squares is:
$\left(\Delta x_{1}\right)^{2}=\left(x_{1}-x\right)^{2}$
$\left(\Delta x_{2}\right)^{2}=\left(x_{2}-x\right)^{2}$
$\left(\Delta x_{3}\right)^{2}=\left(x_{3}-x\right)^{2}$
$\left(\Delta x_{n}\right)^{2}=\left(x_{n}-x\right)^{2}$
If $x$ representes the average of the squares, then the sum of the squares of the errors becomes minimal: $S(x)=\left(x_{1}-x\right)^{2}+\left(x_{2}-x\right)^{2}+\left(x_{3}-x\right)^{2}+\ldots\left(x_{n}-x\right)^{2}$

The derivative with respect to x is canceled and it can be obtained as follows:
$x=\bar{x}=\frac{\sum_{i=1}^{n} x_{i}}{n}$ which is called THE AVERAGE VALUE

It can be determined as follows: THE MEAN SQUARED ERROR :

$$
m= \pm \sqrt{\frac{\left(\Delta x_{1}\right)^{2}+\left(\Delta x_{2}\right)^{2}+\left(\Delta x_{3}\right)^{2}+\ldots .+\left(\Delta x_{n}\right)^{2}}{n}}
$$

or $m= \pm \sqrt{\frac{\sum_{i=1}^{n}\left(\Delta x_{i}\right)^{2}}{n}}$

This parameter depends very little on the number of measurements, if a sufficient amount of determinations are made.
$\mathbf{m}^{\mathbf{2}}$ is called DISPERSION - it represents the square of the mean squared error.

## APPARENT ERRORS :

```
\(v_{1}=x_{1}-\bar{x}\)
\(v_{2}=x_{2}-\bar{x}\)
\(\nu_{3}=x_{3}-\bar{x}\)
\(v_{n}=x_{n}-\bar{x}\)
---------
\(\sum_{i=1}^{n} v_{i}=\sum_{i=1}^{n} x_{i}-n \bar{x}=\sum x_{i}-\frac{n \sum x_{i}}{n}\)
\(\sum_{i=1}^{n} v_{i}=0\)
```

The link between the real and the apparent error is the following:
$\sum x_{i}-v_{i}=\bar{x}-n=\varepsilon$
The mean squared error can also be expressed by using the apparent errors: $\sum x_{i}=v_{i}+\varepsilon$

$$
\begin{aligned}
& \varepsilon=\Delta x_{1}-v_{1} \\
& \varepsilon=\Delta x_{2}-v_{2} \\
& \ldots \ldots \ldots . . . . . \\
& \varepsilon=\Delta x_{n}-v_{n} \\
& \text { nع }=\sum \Delta x_{i}-\sum v_{i}=\sum \Delta x_{i} \\
& \text { so } \varepsilon=\frac{\sum \Delta x_{i}}{n}
\end{aligned}
$$

$$
\sum\left(\Delta x_{i}\right)^{2}=\sum v_{i}^{2}+2 \varepsilon \sum v_{i}+n \varepsilon^{2}
$$

$$
\sum\left(\Delta x_{i}\right)^{2}=\sum v_{i}^{2}+n \varepsilon^{2}
$$

$$
\frac{\sum\left(\Delta x_{i}\right)^{2}}{n}=\frac{\sum v_{i}^{2}}{n}+\frac{\sum\left(\Delta x_{i}^{2}\right)}{n}=\frac{\sum v_{i}^{2}}{n}+\frac{\sum\left(\Delta x_{i}^{2}\right)}{n}+\frac{2 \sum_{i} \sum_{j} \Delta x_{i} \Delta x_{j}}{n^{2}}
$$

as a result, the folowing formula is obtained: $m^{2}=\frac{\sum\left(\Delta x_{i}\right)^{2}}{n}$
afterwards, the following formula is obtained for the mean squared error: $m= \pm \sqrt{\frac{\sum v_{i}^{2}}{n-1}}$
where $v_{i}=x_{i}-\bar{x}$ and $\boldsymbol{n}$ represents the number of measurements.
THE MEAN SQUARED ERROR becomes: $m= \pm \sqrt{\frac{\sum\left(x_{i}-\bar{x}\right)^{2}}{n-1}}$
If a certain value of a measurement xi and $m$ is known, then it is expected that the real value is included in the interval: $\left(\mathrm{x}_{\mathrm{i}}-\mathrm{m}, \mathrm{x}_{\mathrm{i}}+\mathrm{m}\right)$.

Because the average value does not coincide with the real value, it presents a characteristic deviation from it, also called the MEAN SQUARED ERROR OF AVERAGE: $\mu= \pm \sqrt{\frac{\sum\left(x_{i}-\bar{x}\right)^{2}}{n(n-1)}}$.

The actual value will be included in the following interval $(\bar{x}-\mu, \bar{x}+\mu)$.

The properties of accidental errors can be expressed using the Gaussian distribution function, if a continuous function will be accepted for accidental errors:

$$
P(v)=\frac{1}{\sigma \sqrt{2 \pi .}} \exp \left(-\frac{v^{2}}{2 \sigma^{2}}\right)
$$

where $\sigma$ is a constant that depends on the nature of the measurements and the precision of the executions
and $\mathrm{P}(\mathrm{v})$ represents the probability density or the probability that the error of a quantity is included in the interval ( $\mathrm{v}, \mathrm{v}+\mathrm{dv}$ ).

The probability that the error is included in the interval $(-\delta,+\delta)$ is:
$\alpha=\int_{-\delta}^{+\delta} P(v) d v=\frac{1}{\sigma \sqrt{2 \pi}} \int_{-\delta}^{+\delta} \exp \left(-\frac{v^{2}}{2 \sigma^{2}}\right) d v=$
$=\frac{1}{\sqrt{2 \pi}} \int_{-\eta}^{+\eta} \exp \left(-\frac{t^{2}}{2}\right) d t=\frac{2}{\sqrt{2 \pi}} \int_{0}^{+\eta} \exp \left(-\frac{t^{2}}{2}\right) d t$

Where $\eta=\frac{\delta}{\sigma} \quad ; \quad v=\sigma . t$
Performing the integral we arrive at a sum calculated in the tables for different values of $\eta$.
The interpretation of $\sigma$ is as follows: if $\eta=1$, then $\sigma=\delta$.
Thus, $2 \sigma$ represents the interval $(-\sigma,+\sigma)$, in which approximately $68 \%$ of the errors are included, that represents a fraction of $2 / 3$ of the number of errors.
$\alpha$ - represents the probability that the error is included in the interval $(-\delta,+\delta)$;
$\alpha$ - is called confidence coefficient.
$(-\delta,+\delta)$ - confidence interval.
Thus, the probability that the real value x is included in the interval $\left(\mathrm{x}_{\mathrm{i}}-\mathrm{m}, \mathrm{x}_{\mathrm{i}}+\mathrm{m}\right)$ represents 68\%.

The mean squared error of average has the following interpretation:

$$
\mu= \pm \sqrt{\frac{\sum\left(x_{i}-\bar{x}\right)^{2}}{n(n-1)}}: \text { there is a probability of } 68 \%(\approx 2 / 3) \text { for the real value of the measured }
$$ quantity to be included in the interval $(\bar{x}-\mu, \bar{x}+\mu)$ and the results are expressed in the form: $x=\bar{x} \pm \mu$

Sometimes the aim is to find the interval $(-\delta,+\delta)$ in which $50 \%$ of the errors are included.
In this case, the value obtained for a single measurement is the mean squared error: $m= \pm \frac{2}{3} \sqrt{\frac{\sum v_{i}^{2}}{n-1}}= \pm \frac{2}{3} \sqrt{\frac{\sum\left(x_{i}-\bar{x}\right)^{2}}{n-1}}$ and the mean squared error of average: $\mu_{p}= \pm \frac{2}{3} \sqrt{\frac{\sum\left(x_{i}-\bar{x}\right)^{2}}{n(n-1)}}$ which are called probable errors.

In order to obtain a clear overview of the accidental errors, it is useful to use the relative errors defined for each individual measurement by the ratios: $\mathrm{v}_{1} / \mathrm{x}_{1}, \mathrm{v}_{2} / \mathrm{x}_{2}, \ldots \ldots, \mathrm{v}_{\mathrm{n}} / \mathrm{x}_{\mathrm{n}}$. And for the arithmetic mean $\pm \frac{m}{x}$.

If these two types of relationships are multiplied by 100, we get the relative percentage values.

## 2. CALCULATION OF ERRORS FOR QUANTITIES RESULTING FROM A CALCULATION

In most cases, the aim is to determine a quantity obtained on the basis of direct measurements. Knowing the errors committed when measuring the quantities that may intervene in obtaining the final result, a method should be found to calculate the error for the result itself. Since the errors are very small, the serial development method can be used.

The theory of errors states that if a large number of measurements are made directly on each quantity that enters a formula, the mean squared error of the averege of the quantity as a result of a calculation is given by the relationship $f\left(\mathrm{x}_{1}, \mathrm{x}_{2}, \mathrm{x}_{3}, \ldots, \mathrm{x}_{\mathrm{n}}\right)=\mathrm{Y}$ the quantity resulting from a calculation made with the help of directly measurable quantities $\mathrm{x}_{1}, \mathrm{x}_{2}, \ldots, \mathrm{x}_{\mathrm{n}}$.

And :

$$
\mu= \pm \sqrt{\left(\frac{\partial f}{\partial x_{1}} \mu_{1}\right)^{2}+\left(\frac{\partial f}{\partial x_{2}} \mu_{n 1}\right)^{2}+\ldots+\left(\frac{\partial f}{\partial x_{n}} \mu_{n}\right)^{2}}
$$

it is observed that the relation: $m= \pm \sqrt{\frac{\sum v_{i}^{2}}{n(n-1)}}$ is a particular case of the above formula for $\mathrm{f}=\mathrm{x}_{\mathrm{i}}=\mathrm{x}$ and $\mu_{1}^{2}=\frac{\sum v_{i 1}^{2}}{n_{1}\left(n_{1}-1\right)} ; \quad \mu_{2}^{2}=\frac{\sum v_{i 2}^{2}}{n_{2}\left(n_{2}-1\right)} ; \quad \ldots . \quad ; \quad \mu_{N}^{2}=\frac{\sum v_{i N}^{2}}{n_{N}\left(n_{N}-1\right)}$

## 3. LEAST SQUARES METHOD

There are laboratory works in which the aim is to determine a quantity that constitutes a constant coefficient from a linear relationship of the form: $\mathbf{Y}=\mathbf{A} \cdot \mathbf{x}+\mathbf{B}$.

Such a relation is Hooke's law, for example: $F=k x$, in which the determination of the coefficient $k$ or the acceleration of the center of the cylinder rolling on an inclined plane with rolling friction is required. Depending on $\sin \alpha \approx \alpha($ for $\alpha-\operatorname{small}: \mathrm{a}=\mathrm{A} . \alpha+\mathrm{B})$ where B is required.

By making several measurements for $x$ and $y$, the constants $A$ and $B$ can be determined from a graphic representation. Usually the representative points are not exactly on a straight line, being affected by accidental errors. The question arises by which methods the constants A and B should be determined so that their values are as close as possible to the real value.

This problem is solved by the "least squares method" which is based on the following condition: the most probable value of the constants will be the one for which the sum of the squared errors is minimal.

It can be written accordingly:

$$
\begin{aligned}
& \bar{x}=\frac{1}{n} \sum_{i} x_{i} \\
& \bar{y}=\frac{1}{n} \sum_{i} y_{i} \\
& \overline{x^{2}}=\frac{1}{n} \sum_{i} x_{i}^{2} \\
& \overline{x y}=\frac{1}{n} \sum_{i} x_{i} y_{i} \\
& \text { So: } \begin{array}{l}
\overline{x^{2}} A+\bar{x} B-\overline{x y}=0 \\
\bar{x} A+B-\bar{y}=0
\end{array}
\end{aligned}
$$

The solutions of these equations constitute the values sought for A and B :

$$
\begin{aligned}
& A=\frac{\overline{x y}-\bar{x} \cdot \bar{y}}{\overline{x^{2}}-(\bar{x})^{2}} \\
& B=\frac{\overline{x^{2}} \cdot \bar{y}-\bar{x} \cdot \overline{x y}}{\overline{x^{2}}-(\bar{x})^{2}}
\end{aligned}
$$

The method can also be applied in non-linear relationships obtained on the basis of similar demonstrations.
4. STATISTICAL METHOD OF CALCULATION FOR THE AVERAGE VALUE $\overline{\mathcal{X}}$

1. The statistical calculation mode is entered

$$
\text { Mode } \rightarrow \text { STAT }
$$

2. Enter the data into the statistics:

3. The number of data entered into the statistic should be checked: $n=5$

$$
(\mathrm{SHIFT})(\mathrm{n}) \Rightarrow \mathrm{n}=5
$$

4. The average value is calculated:

$$
\bar{x}=10.51
$$

## Statistical method of calculation for error

$\mu= \pm \sqrt{\frac{\sum_{i=1}^{n}\left(x_{i}-\bar{x}\right)^{2}}{n(n-1)}}$
$\mu= \pm \sqrt{\frac{(10,52-10,51)^{2}+(10,65-10,51)^{2}+(10,45-10,51)^{2}+(10,45-10,51)^{2}+(10,48-10,51)^{2}}{5 \cdot 4}}$

1. The statistical calculation mode is entered
2. The previous memory should be deleted
3. Enter the differences in memory:

$$
\begin{aligned}
& (10.52-10.51)=\text { DATE } \\
& (10.65-10.51)=\text { DATE } \\
& (10.45-10.51)=\text { DATE } \\
& (10.45-10.51)=\text { DATE }
\end{aligned}
$$

$$
(10.48-10.51)=\text { DATE }
$$

4. Check no. of data from memory: $\mathrm{n}=5$

$$
(\mathrm{SHIFT})(6) \Rightarrow \mathrm{n}=5
$$

5. The sum of squares is calculated

$$
\sum x^{2}=0,0278
$$

6. Afterwards, the following calculation is made $\frac{\sum x^{2}}{5 \cdot 4}=1,39 \cdot 10^{-3}$
7. It is calculated $\sqrt{\frac{\sum x^{2}}{5 \cdot 4}}=\sqrt{1,39 \cdot 10^{-3}}=0,0372$
8. The calculated error is $\mu=0.0372$

Thus, the range of actual values is:

$$
\begin{aligned}
& x_{\text {real }} \in(10,51-0,0372 ; 10,51+0,0372) \\
& x_{\text {real }} \in(10,4728 ; 10,5472) m
\end{aligned}
$$

## DETERMINATION OF DENSITY FOR SOLID MATERIALS

## INTRODUCTION

The density of a material is an old physical notion since the century. XVII. Émile Littré defined density as the ratio between mass and volume. Laplace uses in 1810 for the notion of density it of volumetric mass, which is still associated in 1960 by the International Bureau of Weights and Measures.

Density is the ratio between mass and its volume.
The relative density of a material is the ratio between its density and the density of a reference material. In the case of liquids, pure water at $4{ }^{\circ} \mathrm{C}$ is used as a reference. In the case of liquids, the precise measurement of the density is carried out with the help of the pycnometer. In the case of measuring the density of gases or vapors, the reference gas is air. Density is a quantity that has the following unit of measure: $\mathrm{kg} / \mathrm{m} 3$, and relative density has no unit of measure.

Densimeters are used for research and development, quality control, commercial transactions, in industrial sectors such as:

- petrochemicals : numerous rules define the use of different models of densimeters in this sector.
- Pharmacy : according to the methods described in the European Pharmacopoeia.
- Beverages: to determine the concentration of sugar of sweet drinks, and for manufacture beer, to determine the level of alcohol in alcoholic beverages;
- Chemistry: by measuring the density, numerous chemical products are allowed such as: $\mathrm{H}_{2}$ $\mathrm{SO}_{4}, \mathrm{NaOH}, \mathrm{HNO}_{3}$, etc. ;
- Cosmetics: the production of cosmetic items requires numerous checks on raw materials and finished products in compliance with the laws


## THE DENSITY

It is a physical quantity that characterizes the mass of the volume unit. It is generally noted with the Greek letter $\rho$ ( $r h \hat{o}$ ) as it is also used by the International Bureau of Weights and Measures.

Density is by definition the ratio between mass and volume $\rho=\frac{m}{V}$
Where:
m - the mass of the studied material
V- the volume of the studied material

Density is synonymous with absolute density and specific density

The unit of measurement in the International System is: kilogram per cubic meter: ( kg $\cdot \mathrm{m}^{-3}$ or $\left.\mathrm{kg} / \mathrm{m}^{3}\right):[\rho]_{S I}=\frac{k g}{m^{3}}$

The unit of measurement in the CGS System is: gram per cubic centimeter: ( $\mathrm{g} \cdot \mathrm{cm}^{-}$
${ }^{3}$ or g/cm $\left.{ }^{3}\right):[\rho]_{C G S}=\frac{g}{\mathrm{~cm}^{3}}$
Apparent density or bulk density is defined as the ratio of the mass of the material and the volume of the grain assembly: $\rho_{\text {apparent }}=\frac{m_{\text {mat }}}{V_{\text {apparent }}}$

The real density is the ratio between the mass of the material and the real mass of the grains

$$
\rho_{\text {real }}=\frac{m_{\text {mat }}}{V_{\text {real }}}
$$

## Absolute density - is used for porous materials.

This size is used in the case of porous materials. For these, the material must be finely ground and the actual density of the powder obtained is measured. The absolute density of matter is the ratio between the mass of the material and the real volume for which the volume of the pores (open and closed) has been subtracted.

Absolute density is equal to particle density for non-porous materials.

$$
\rho_{\text {absolute }}=\frac{m_{\text {mat }}}{V_{\text {absolute }}}=\frac{m_{\text {mat }}}{V_{\text {real }}-V_{\text {pore }}}
$$

Table 1: various values of the density of solid materials

| Mineral rocks common materials | density $\left(\mathbf{k g} / \mathbf{m}^{\mathbf{3}}\right)$ |
| :--- | :--- |
| asbestos | 2500 |
| clay | $1300-1700$ |
| limestone | $2000-2800$ |
| compost | $550-600$ |
| chalk | $1700-2100$ |
| diamond | 3517 |
| granite | $1800-2500$ |
| floor tiles | $1600-1900$ |
| kaolin | 2260 |
| marble | $2650-2750$ |
| cjart | 2650 |
| pumice | 910 |
| porcelain | 2500 |
| sand | $1600(\mathrm{sec})-2000($ saturated $)$ |
| silicon | 2330 |
| vegetable soil | 1250 |
| glass range | 2530 |
| cotton | $20-60$ |
| Metals and alloys | density $\left(\mathbf{k g} / \mathbf{m}^{\mathbf{3}}\right)$ |
| steel | $7,500-8,100$ |
| high speed steel | $8,400-9,000$ |
| cast iron | $6,800-7,400$ |
| aluminum | 2700 |
| silver | 10,500 |
| beryllium | 1848 |
| bronze | $8,400-9,200$ |
| Carbon (diamond) | 3508 |
| carbon (graphite) | 2250 |
| constantly | 8910 |
| copper | 8920 |
| duralumin | 2900 |
| tin | 290 |
|  |  |


| Mineral rocks common materials | density $\left(\mathbf{k g} / \mathbf{m}^{\mathbf{3}}\right)$ |
| :--- | :--- |
| iron | 7,860 |
| iridium | 22560 |
| brass | $7300-8800$ |
| lithium | 530 |
| magnesium | 1750 |
| Mercury | 13545.88 |
| molybdenum | 10,200 |
| nickel | 8,900 |
| gold | 19,300 |
| osmium | 22610 |
| palladium | 12,000 |
| platinum | 21450 |
| lead | 11350 |
| potassium | 850 |
| tantalus | 16,600 |
| titanium | 4500 |
| tungsten | 19,300 |
| uranium | 18,700 |
| vanadium | 6100 |
| zinc | 150 |
|  |  |

## MEASUREMENTS WITH BALANCE

The operating principle of electronic balances is based on the variations in the electrical characteristics of some materials when they are subjected to mechanical compression. The electrical variations are measured and recorded on a screen that allows the expression of the mass of the unknown body that compresses the component elements of the balance.

Depending on the masses that can be measured, electronic balances are classified according to divisions and measuring capacity:

Example of analytical balances:

balance ( $0.1 \mu \mathrm{~g} / 3 \mathrm{~g}$ );
balance ( $0.001 \mathrm{mg} / 3 \mathrm{~g}$ );
balance ( $0.01 \mathrm{mg} / 30 \mathrm{~g}$ );
balance ( $0.1 \mathrm{mg} / 160 \mathrm{~g}$ );
precision balance ( $1 \mathrm{mg} / 160 \mathrm{~g}-60 \mathrm{~kg}$ )

## Hand balance



It is used to weigh small amounts of substance: usually $0.06 \mathrm{~g}-100-200 \mathrm{~g}$, several weighing models are known of 0.25-0.5-1-2-5-10-20 g (with the sensitivity decreasing).

Constructively, the hand balance consists of a suspended metal lever, with the arrow pointing upwards, with two equal arms that oscillate around a support point finished with a ring.

At the ends of the arms, 2 anti-corrosive plates, usually made of celluloid or plastic material, are suspended in 3 silk cords.
Hand balance can be simple or with a slider.
With hand balances, weighing is done by balancing the weight of the substance on one of the plates with corresponding weights placed on the opposite platen. On the hand balance with a slider, one of the arms, usually the left one, is divided, the cursor sliding along the lever that indicates the weight of the substance to be weighed (placed on the platen corresponding to the undivided arm.

Sometimes additional weights can be used on the plate of the graduated arm. This operation is not recommended especially for weighing strongly active substances.

## Pharmaceutical balance:

## Or balance with column:



It is used in the pharmacy to weigh from 10 to 1000 g or from 5 to 500 g with a sensitivity of centigrams.

The balance consists of a metal column, mounted vertically on a box or support, usually covered with a marble plate, with 1-3 drawers. On the upper part of the column is placed on a knife, a horizontal lever with equal arms, on which two plates are suspended on the knives.

The process of balancing of the instrument is signaled by the indicator needle, positioned in the middle of the lever and which oscillates together with it in front of a graduated dial. The process of balancing of the instrument is checked before each weighing, the pharmaceutical balance having adjustment screws at the end of each arm.


## Weighing with a balance with unequal arms

$$
\begin{aligned}
& m \cdot g \cdot r_{a}=m_{1} \cdot g \cdot r_{b} \\
& m_{2} \cdot g \cdot r_{a}=m \cdot g \cdot r_{b} \\
& \frac{m}{m_{2}}=\frac{m_{1}}{m} \\
& m=\sqrt{m_{1} \cdot m_{2}}
\end{aligned}
$$

## MEASURING USING VERNIER CALIPERS

The caliper is an instrument for measuring lengths whose construction is based on the principle of the linear vernier. The precision of the measurements made with the caliper varies between $\frac{1}{10}$ and $\frac{1}{50} \mathrm{~mm}$.

## Instrument description

From the adjacent figure it can be seen that the caliper is made up of the following parts:


- a metal ruler $R$, graduated in $m m$, which ends with a stop device $A$ on one side and another stop device $A^{\prime}$, on the opposite side;
- a movable cylinder head $C$;
- a vernier $V$ applied to the cylinder head;
- and cylinder head $B^{\prime}$ devices $B$ that correspond to the stop devices $A$ and $A^{\prime}$ that serve to delimit the dimensions to be measured;
- movable blade $P$, attached to the cylinder head, the blade that moves with the vernier and that serves to measure depths, while it $A B$ serves to delimit thicknesses, external diameters, etc., and $A^{\prime} B^{\prime}$ serves to determine internal diameters;
- a screw $S$ for immobilizing the cylinder head.
a) Measuring the inner diameter for cylindrical parts:

It is done in such a way that the stop device is fixed inside the cylinder $A^{\prime}$ at one end of the inner diameter of the cylinder, and by moving the cylinder head, the mobile device is fixed $B^{\prime}$ at the other end of the inner diameter. The cylinder head screw is screwed in $S$ and then the reading is taken.
b) Measurement of the outer diameter for cylindrical parts :

The measurement of the external diameter for cylindrical parts of the thickness of a bar, or the thickness of a metal plate, etc., is done as follows: the cylinder head of the vernier is moved and the part to be measured is clamped between the branches and so that the part is well fixed, without deforming $A$ it $B$. Fix the cylinder head bolt so that it does not move and then take the reading according to the indications given as follows.

## c) Measuring the depth of a piece:

It is done as follows: to accurately measure the depth of the cylinder, the cylinder head of the vernier is moved. The blade moves together with the cylinder head. When the base of the cylinder has been reached, resting the upper edge on the ruler of the caliper, tighten the screw $S$ and read the dimensions on the ruler and on the vernier.

## PROCEDURE FOR THE PRACTICAL WORK - METHOD 1:

1. The mass of the material to be studied is determined using the analytical balance from the laboratory
2. Calculate the volume of the solid material considered: a ring with inner diameter Di, external diameter De and height h , using the following formula: $\quad V=\frac{\pi}{4} . h .\left(D_{e}^{2}-D_{i}^{2}\right)$

Di, De and h are determined using the caliper.

3. The density is calculated with the formula: $\rho=\frac{m}{V}$.
4.5 measurements are made for $\mathrm{h}, \mathrm{De}, \mathrm{Di}$.
5. 5 calculations are made for V and $\rho$.

6 . Write the values in the table below.
7. The errors and the interval for real values should be calculated.

The average error and the range of real values are calculated
Average value $\rho=\bar{\rho}=\frac{\sum_{i=1}^{n} \rho_{i}}{n}$
THE MEAN SQUARED ERROR OF THE AVERAGE $\mu= \pm \sqrt{\frac{\sum\left(\rho_{i}-\bar{\rho}\right)^{2}}{n(n-1)}}$
Interval for real values: $(\bar{\rho}-\mu, \bar{\rho}+\mu)$

Table 1: determination of the density at the solid ring

| No <br> crt | $\mathrm{h}(\mathrm{cm})$ | $\mathrm{D}_{\text {ext }}$ <br> $(\mathrm{cm})$ | $\mathrm{D}_{\text {int }}$ <br> $(\mathrm{cm})$ | V <br> $(\mathrm{cm} 3)$ | m <br> $(\mathrm{g})$ | $\rho\left(\frac{g}{\mathrm{~cm}^{3}}\right)$ | $\bar{\rho}\left(\frac{g}{\mathrm{~cm}^{3}}\right)$ | $\boldsymbol{\mu}$ | $\rho_{\text {real }}$ |
| :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- |
| $\mathbf{1}$ |  |  |  |  |  |  |  |  |  |
| $\mathbf{2}$ |  |  |  |  |  |  |  |  |  |
| $\mathbf{3}$ |  |  |  |  |  |  |  |  |  |
| $\mathbf{4}$ |  |  |  |  |  |  |  |  |  |
| $\mathbf{5}$ |  |  |  |  |  |  |  |  |  |

PROCEDURE FOR THE PRACTICAL WORK - METHOD 2:
LECITHIN, ESENTIN, AMINO ACID Pills or others may be choosen

1. The mass of the studied material is determined with the help of the balance.
2. The diameter and height are measured using a caliper and based on these values we can calculate the volume of the lecithin tablets that is formed by a cylinder +2 semispheres at the ends :
$\mathrm{V}=\frac{\pi \cdot \mathrm{D}^{2}}{4}(\mathrm{~h}-\mathrm{D})+\frac{\pi \cdot \mathrm{D}^{3}}{6}$
3. With the obtained volume and masses we are able to calculate the density, using the following formula : $\rho=\frac{m}{v}$.

4. Make 5 measurements for h and D .
5. 5 calculations should be made for each V and $\rho$.
6. The average density value should be calculated using the arithmetic mean formula $\bar{\rho}=\frac{\sum_{i=1}^{n} \rho_{i}}{n}$
7. The mean squared error of the average should be calculated using the following formula: $\mu= \pm \sqrt{\frac{\sum\left(\rho_{i}-\bar{\rho}\right)^{2}}{n(n-1)}}$
8. Calculate the interval of real values $(\bar{\rho}-\mu, \bar{\rho}+\mu)$.

Table 2: determination of the density of lecithin tablets

| No. <br> crt. | h | d | $\mathrm{V}(\mathrm{cm} 3)$ | $\mathrm{m}(\mathrm{g})$ | $\rho\left(\frac{g}{\mathrm{~cm}^{3}}\right)$ | $\bar{\rho}\left(\frac{g}{\mathrm{~cm}^{3}}\right)$ | $\boldsymbol{\mu}$ | $\rho_{\text {real }}$ |
| :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- |
| $\mathbf{1}$ |  |  |  |  |  |  |  |  |
| $\mathbf{2}$ |  |  |  |  |  |  |  |  |
| $\mathbf{3}$ |  |  |  |  |  |  |  |  |
| $\mathbf{4}$ |  |  |  |  |  |  |  |  |
| $\mathbf{5}$ |  |  |  |  |  |  |  |  |
|  |  |  |  |  |  |  |  |  |

CONCLUSIONS Errors for directly measured quantities can be calculated. Also, the errors of the quantities obtained from a calculation can be determined and the results can be compared. Similar calculations can be made for materials with different shapes.

## DETERMINATION OF LIQUID DENSITY USING THE PYCNOMETER

## INTRODUCTION

Density (after an older name, "specific gravity") is an intensive property of substances. By definition, absolute density is the mass of the unit volume of the respective substance. If a material (solid, liquid or gas) with the mass " $m$ " occupies the volume " V ", then the absolute density (" $\rho$ ") of the body is:

$$
\rho=\frac{\mathrm{m}}{\mathrm{~V}}(1)
$$

In the International System of Units (SI), absolute density is expressed in $\mathrm{kg} / \mathrm{m}^{3}$. In practice, tolerated measurement units are often used ( $\mathrm{g} / \mathrm{cm}^{3}, \mathrm{~g} / \mathrm{liter}$, etc.).

Due to the presence of the volume " V " in the definition relation (1), the value of the absolute density of a material (solid or liquid) depends on the temperature.

The relative density " d " is often used, and it represents the ratio between the absolute density of the substance of interest $(\rho)$ and the absolute density of a reference material ( $\rho_{\text {ref }}$ ).

$$
d=\frac{\rho}{\rho_{\text {ref }}}(2)
$$

From relation (2) results that the relative density is a adimensional quantity. Since both quantities depend on the temperature, the relative density also depends on the temperature. Water density is frequently chosen as a " $\rho_{\text {ref }}$ ". Regarding the temperature, two expressions of this quantity are used. If one compares the absolute density of the studied medium (" $\rho$ "), measured at a temperature of $20^{\circ} \mathrm{C}$, with the absolute density (" $\rho_{\text {ref }}$ ") of the reference medium, also measured at a temperature of $20^{\circ} \mathrm{C}$, then the relative density is noted as " $d_{20}^{20}$ ". If you compare the absolute density of the studied medium (" $\rho$ "), measured at a temperature of $20^{\circ} \mathrm{C}$, with the absolute density ( $" \rho_{\text {ref }}$ ") of water, measured at a temperature of $4^{\circ} \mathrm{C}$, then the relative density is noted as " $d_{4}^{20}$ ". The relevance of this latter comparison lies in the fact that the density of water is maximum at a temperature of $4{ }^{\circ} \mathrm{C}$.

Determination of the density of a liquid by means of the pycnometer

Figure 2 represents a simple pycnometer, and Figure 3 represents a pycnometer equipped with a thermometer. The pycnometer shown in Figure 2 consists of a glass flask ("b") and a glass stopper ("d") provided with a thin channel ("c").

(A)

(B)

(A)

(B)

Fig. 2 Simple pycnometer Fig. 3 Pycnometer equipped with thermometer

## PROCEDURE FOR THE PRACTICAL WORK - STEPS

1) Weigh the empty and dry pycnometer, together with the stopper "d" (Figure 2-(A), mass " $m_{0}$ ").
2) A reference liquid for which the absolute density " $\rho_{\text {ref }}$ " is known is introduced into the balloon "b", so that the balloon " b " is filled to the brim, after which it is closed with the stopper " d " (Figure 2B ). The excess of reference liquid comes out through channel "c", so the volume of reference liquid remaining in the pycnometer is equal to the sum of the volume of balloon "b" and the volume of channel "c" (volume "V").
3) Weigh the pycnometer filled with reference liquid (" $m_{\text {ref }}$ ").
4) Empty, wash and dry the pycnometer, after which it should be filled with the studied liquid in the manner described at step 2.
5) Weigh the flask filled with the studied liquid ("m").
6) The quantities involved satisfy relation (5).

$$
\mathrm{V}=\frac{\mathrm{m}_{\mathrm{ref}}-\mathrm{m}_{\mathrm{o}}}{\rho_{\mathrm{ref}}}=\frac{\mathrm{m}-\mathrm{m}_{\mathrm{o}}}{\rho}(5)
$$

$\rho$ - the absolute density of the studied liquid can be expressed as follows :

$$
\rho=\rho_{\mathrm{ref}} \cdot \frac{\mathrm{~m}-\mathrm{m}_{\mathrm{o}}}{\mathrm{~m}_{\mathrm{ref}}-\mathrm{m}_{\mathrm{o}}}(6)
$$

and also the relative density " d ":

$$
\mathrm{d}=\frac{\rho}{\rho_{\mathrm{ref}}}=\frac{\mathrm{m}-\mathrm{m}_{\mathrm{o}}}{\mathrm{~m}_{\mathrm{ref}}-\mathrm{m}_{\mathrm{o}}}(7)
$$

The experimental part is performed by using three liquids, ethanol $\left(\mathrm{C}_{2} \mathrm{H}_{5}-\mathrm{OH}\right)$, carbon tetrachloride $\left(\mathrm{CCl}_{4}\right)$ and acetone $\left(\left(\mathrm{CH}_{3}\right)_{2} \mathrm{CO}\right)$, the reference liquid being distilled water each time. The density of water as a function of temperature is shown in Table 2 and Figure 4.

Table 2 Density of water depending on temperature

| $\mathbf{t}\left({ }^{\mathbf{}} \mathbf{C}\right)$ | $\boldsymbol{\rho}\left(\mathbf{k g} / \mathbf{m}^{\mathbf{3}}\right)$ | $\mathbf{t}\left({ }^{( } \mathbf{C}\right)$ | $\boldsymbol{\rho}\left(\mathbf{k g} / \mathbf{m}^{\mathbf{3}}\right)$ | $\mathbf{t}\left({ }^{( } \mathbf{C}\right)$ | $\boldsymbol{\rho}\left(\mathbf{k g} / \mathbf{m}^{\mathbf{3}}\right)$ |
| :--- | :--- | :--- | :--- | :--- | :--- |
| $\mathbf{0}$ | 999.87 | $\mathbf{1 3}$ | 999.40 | $\mathbf{2 2}$ | 997.80 |
| $\mathbf{1}$ | 999.93 | $\mathbf{1 4}$ | 999.27 | $\mathbf{2 . 3}$ | 997.56 |
| $\mathbf{2}$ | 999.97 | $\mathbf{1 5}$ | 999.13 | $\mathbf{2 4}$ | 997.32 |
| $\mathbf{3}$ | 999.99 | $\mathbf{1 6}$ | 998.97 | $\mathbf{2 5}$ | 997.07 |
| $\mathbf{4}$ | 1000.00 | $\mathbf{1 7}$ | 998.80 | $\mathbf{2 6}$ | 996.81 |
| $\mathbf{5}$ | 999.99 | $\mathbf{1 8}$ | 998.62 | $\mathbf{2 7}$ | 996.54 |
| $\mathbf{6}$ | 999.97 | $\mathbf{1 9}$ | 998.43 | $\mathbf{2 8}$ | 996.26 |
| $\mathbf{7}$ | 999.93 | $\mathbf{2 0}$ | 998.23 | $\mathbf{2 9}$ | 995.97 |
| $\mathbf{8}$ | 999.88 | $\mathbf{2 1}$ | 998.02 | $\mathbf{3 0}$ | 995.67 |



Fig. 4 Density of water depending on temperature

The data obtained form the experimental part should be introduced in Table 3.1 and Table 3.2.

Table 3.1.

| Liquid | $\mathbf{m}(\mathbf{g})$ | $\rho\left(\mathbf{g} / \mathbf{c m}^{\mathbf{3}}\right)$ |
| :--- | :--- | :--- |
| $\mathrm{C}_{2} \mathrm{H}_{5} \mathrm{OH}$ |  |  |
| $\mathrm{CCl}_{4}$ |  |  |
| $\left(\mathrm{CH}_{3}\right)_{2} \mathrm{CO}$ |  |  |

Table 3.2. Determination of liquid density using the pycnometer:
$\mathbf{N a C l}$ solution of different concentrations, $\mathrm{m}_{0}=$ g
$\mathrm{V}=\mathbf{2 5} \mathbf{~ m l}, \mathbf{m}(\mathrm{g})=$ mass of the pycnometer filled with liquid

| $\left.\mathbf{C N a C l}^{(\%)}\right)$ | $\mathbf{M}(\mathrm{g})$ | $\rho\left(\mathbf{g} / \mathrm{cm}^{3}\right)$ |
| :--- | :--- | :--- |
| 0.9 |  |  |
| 1 |  |  |
| 5 |  |  |
| 6 |  |  |
| 8 |  |  |
| 10 |  |  |
| 15 |  |  |
| 20 |  |  |
| 25 |  |  |
| 30 |  |  |
| 50 |  |  |

## DETERMINATION OF THE DENSITY OF A SOLID MATERIAL USING THE PYCNOMETER

If the size of a solid material can be entered inside a pycnometer, then the density of the body can be determined using the pycnometer. This technique is especially useful if the studied sample is a wax or has the consistency of a paste, or consists of several particles, not necessarily identical. In this case, an average (apparent) density of the solid or semi-solid sample can be determined.

## PROCEDURE FOR THE PRACTICAL WORK - STEPS

1) Weigh the empty dry pycnometer (together with the stopper, Figure 5-1; mass " $m_{0}$ ").
2) Fill the pycnometer with a reference liquid for which the density " $\rho_{\mathrm{r}}$ " is known (Figure 5-2; mass " $\mathrm{m}_{\mathrm{r}}$ ").
3) Empty, wash and dry the pycnometer, insert the sample into the pycnometer and weigh it (together with the stopper, Figure 5-3; mass " $\mathrm{m}_{\mathrm{p}}$ ")
4) Keeping the sample solid in the pycnometer, fill the volume of the pycnometer with the reference liquid taking care not to form air bubbles and close the pycnometer with the stopper (Figure 5-4, mass " $\mathrm{mpr}^{\mathrm{r}}$ ).

(2) $\left(m_{1}\right)$

(3) $m_{p}$ )


Fig. 5 The stages of determining the density of a solid body with a pycnometer.

- " $\rho$ " solid body density;
$-" \rho_{r}$ " the density of the reference liquid;
- "V" the internal volume (including the volume of the channel in the plug) of the empty pycnometer;
- " $\mathrm{V}_{\mathrm{p}}$ " the volume of the solid sample;
- " $\mathrm{V}_{\mathrm{L}}$ " the volume of reference liquid that fills the pycnometer containing the solid body;
- " $\mathrm{m}_{\mathrm{o}}$ ", " $\mathrm{m}_{\mathrm{r}}$ ", " $\mathrm{m}_{\mathrm{p}}$ " and " $\mathrm{m}_{\mathrm{pr}}$ " the masses weighed in the four stages,
then from the balance of masses and volumes the following relationships result:
$\mathrm{V}_{\mathrm{L}} \cdot \rho_{\mathrm{r}}+\mathrm{V}_{\mathrm{p}} \cdot \rho=\mathrm{m}_{\mathrm{pr}}-\mathrm{m}_{\mathrm{o}}$
$\mathrm{V}_{\mathrm{p}} \cdot \rho=\mathrm{m}_{\mathrm{p}}-\mathrm{m}_{\mathrm{o}}$
Substituting the product $\mathrm{V}_{\mathrm{p}} \cdot \rho$ from relation (9) into relation (8), the following formula results:
$\mathrm{V}_{\mathrm{L}} \cdot \rho_{\mathrm{r}}+\mathrm{m}_{\mathrm{p}}-\mathrm{m}_{\mathrm{o}}=\mathrm{m}_{\mathrm{pr}}-\mathrm{m}_{\mathrm{o}}$
$V_{L}=V-V_{p}=\frac{m_{r}-m_{0}}{\rho_{\mathrm{r}}}-\frac{m_{p}-m_{o}}{\rho}$
On the other hand, $V_{L}=\frac{m_{p r}-m_{p}}{\rho_{r}}$
From relations (11) and (12) the following formula is obtained, by eliminating the volume " $\mathrm{V}_{\mathrm{L}}$ ", $\frac{m_{r}-m_{o}}{\rho_{\mathrm{r}}}-\frac{m_{p}-m_{o}}{\rho}=\frac{m_{p r}-m_{p}}{\rho_{\mathrm{r}}}$

After explaining the density " $\rho$ " from relation (13), we obtain the expression (14) of the density " $\rho$ " of the solid body depending on the weighed masses " $\mathrm{m}_{\mathrm{o}}$ ", " $\mathrm{m}_{\mathrm{r}}$ ", " $\mathrm{m}_{\mathrm{p}}$ " and " $\mathrm{m}_{\mathrm{pr}}$ ": $\rho=\rho_{\mathrm{r}} \cdot \frac{\mathrm{m}_{\mathrm{p}}-\mathrm{m}_{\mathrm{o}}}{\mathrm{m}_{\mathrm{r}}+\mathrm{m}_{\mathrm{p}}-\mathrm{m}_{\mathrm{pr}}-\mathrm{m}_{\mathrm{o}}}$

In practice, the density of a quantity of small glass balls is determined. For this purpose, the operations described in the sequence 1) - 4) are performed. The solid body is an (arbitrary) amount of small glass balls, and distilled water is used as the reference liquid. It is important to check the temperature of the distilled water and use the " $\rho_{\mathrm{r}}$ " value corresponding to that temperature (from Table 2). Weighings are performed with the semimicro analytical balance. The weighed masses and the result are entered in Table 4.

Table 4 Results of solid body density determination

| $\mathrm{m}_{0}(\mathrm{~g})$ | $\mathrm{m}_{\mathrm{r}}(\mathrm{g})$ | $\mathrm{m}_{\mathrm{p}}(\mathrm{g})$ | $\mathrm{mpr}_{\mathrm{pr}}(\mathrm{g})$ | t ( ${ }^{\circ} \mathrm{C}$ ) | $\rho \mathrm{r}\left(\mathrm{g} / \mathrm{cm}^{3}\right)$ | $\rho\left(\mathrm{g} / \mathrm{cm}^{3}\right)$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  |  |  |  |  |  |  |
|  |  |  |  |  |  |  |
|  |  |  |  |  |  |  |

## DETERMINATION OF THE DENSITY OF A LIQUID USING THE MOHR-WESTPHAL BALANCE

## INTRODUCTION

Mohr-Wesrphal balance is also based on Archimedes' principle. Figure 7 shows the construction details of a Mohr-Wesrphal balance.

A float is hung on one of the arms of the balance (in some cases a thermometer is included in the float), and a counterweight is found on the other arm. On the arm equipped with a float there are some recesses where different weights can be fixed.

To determine the density of a liquid, the float is immersed in the reference liquid (for which the density is known; this liquid is often water) and the balance is equilibrated. Lock the position of the balance arms, remove the float, wash and dry. After mounting the float in the initial position, it is immersed in the studied liquid and the balance arms are unlocked. The force of ascent being different from that observed for the reference liquid, the balance becomes unbalanced. By changing the values and positions of the weights on the balance arm, the balance is equilibrated again. Basically, the difference between the apparent weights of the volume of the reference liquid and the column of the studied liquid, displaced by the volume of the float, is determined.


## Fig. 1 Mohr-Westphal balance

If in the first phase of the measurement (when the float is immersed in the reference liquid) the apparent weight of the float is "Gr" and in the second phase (when the float is immersed in the studied liquid) the apparent weight of the float is " $\mathrm{G}_{\mathrm{p}}$ ", then, from the law relating to the ascent force (Archimedes' law) the following formula results:

$$
\mathrm{G}_{\mathrm{p}}-\mathrm{G}_{\mathrm{r}}=\frac{\rho_{\mathrm{p}}-\rho_{\mathrm{r}}}{\rho_{\mathrm{r}}}(1)
$$

In relation (1) " $\rho_{\mathrm{r}}$ " and " $\rho_{\mathrm{p}}$ " respectively represent the density of the reference liquid and the studied liquid (of the sample).

## DENSITY OF A WATER-ETHANOL MIXTURE AS A FUNCTION OF COMPOSITION. VOLUME CONTRACTION WHEN MIXING LIQUIDS

When mixing two miscible liquids, it is found (almost always) that the volume resulting from mixing differs from the sum of the initial volumes (the most common finding is a decrease in the total volume compared to the sum of the mixed volumes - volume contraction). For example, if equal volumes of ethanol and water are mixed, the temperature increases when mixed. After the temperature of the mixture is brought back to the initial value, it is found that the volume of the mixture is smaller than the sum of the mixed volumes.

Figure 2 represents the molar volume of the ethanol-water mixture depending on the composition of the mixture (expressed as a molar fraction of ethanol).


Fig. 2 Molar volume of ethanol-water mixtures according to composition

It is found that for an ethanol-water mixture, in which the molar fraction of ethanol is 0.4 , the molar volume is lower by approximately $1.16 \mathrm{~cm}^{3} / \mathrm{mol}$ than the molar volume resulting from the summation of the volumes. This phenomenon can be followed by determining the density of ethanolwater mixtures depending on the composition of the mixture.

If when mixing the volume " $\mathrm{V}_{1}$ " of ethanol with the volume " $\mathrm{V}_{2}$ " of water, the two liquids having the densities " $\rho_{1}$ " and " $\rho_{2}$ ", then theoretically (if no volume contraction would occur), the density of the mixture would be " $\rho_{\mathrm{t}}$ ", and the volume of the mixture, " $\mathrm{V}_{\mathrm{t}}$ ", would be the sum of the mixed volumes.

$$
\rho_{1} \cdot V_{1}+\rho_{2} \cdot V_{2}=\rho_{t} \cdot\left(V_{1}+V_{2}\right)=\rho_{t} \cdot V_{t}(16)
$$

Practically, due to volume contraction, the volume of the mixture is " $\mathrm{V}_{\mathrm{p}}$ " and the density " $\rho_{\mathrm{p}}$ ".

$$
\rho_{1} \cdot V_{1}+\rho_{2} \cdot V_{2}=\rho_{p} \cdot V_{p}(17)
$$

From relations (16) and (17) it follows:

$$
\rho_{\mathrm{p}}=\rho_{\mathrm{t}} \cdot \frac{\mathrm{~V}_{\mathrm{t}}}{\mathrm{~V}_{\mathrm{p}}}(18)
$$

Since the volume " $V_{p}$ " is lower than the theoretical volume " $V_{t}$ ", the real density " $\rho_{p}$ " is higher than the theoretical density " $\rho \mathrm{t}$ ".

This phenomenon can be highlighted by measuring the real density " $\rho_{p}$ " of the mixture. An alternative way can be based on the use of the semi-micro analytical balance (Figure 3).

## PROCEDURE FOR THE PRACTICAL WORK - STEPS

1) A Berzelius glass, filled with the reference liquid (distilled water), is placed on the pan of an analytical balance and the displayed " $\mathrm{m}_{\text {or }}$ " mass is read (Figure 3-A).
2) A body suspended by a thin wire is immersed in the liquid (Figure 3-B) and the indication on the balance display $\left(\mathrm{m}_{\mathrm{r}}\right)$ is read.
3) Empty, wash and dry both the Berzelius glass and the float, after which the operations from point 1 ) and 2) are repeated, but with the studied liquid in the glass (read the masses $m_{o p}$ and $m_{p}$, respectively (Figures 3-C and 3-D).

The float must be completely immersed in the liquid and not touch the wall or the bottom of the glass.


Fig. 3 Highlighting the volume contraction of the ethanol-water mixture
The experiment is, in a way, the opposite configuration to the one presented on the MohrWesrphal balanc. Since the float is fixed, the force of ascent manifests itself on the glass and implicitly on the pan of the analytical balance.

If the densities of the reference liquid and the studied liquid (sample) are written with " $\rho_{\mathrm{r}}$ " and " $\rho_{\mathrm{p}}$ ", respectively, and with "V" the volume of the float, then, taking into account Archimedes' law, regarding the force of ascent, it can be written: $m_{r}=m_{o r}+V \cdot \rho_{r}$

After removing the volume " V " and explaining the density of the sample " $\rho_{\mathrm{p}}$ ", it results:

$$
\rho_{\mathrm{p}}=\rho_{\mathrm{r}} \cdot \frac{\mathrm{~m}_{\mathrm{p}}-\mathrm{m}_{\mathrm{op}}}{\mathrm{~m}_{\mathrm{r}}-\mathrm{m}_{\mathrm{or}}}(20)
$$

Basically, the density of ethanol-water mixtures is determined, at ambient temperature, for mixtures in which the ethanol:water volumetric ratio is $0: 100$ (pure water); 20:80; 40:60; 60:40; 80:20; 100:0; (pure ethanol). the data are entered in Table 1.

Table 1 Results

| Volumetric ethanol:water ratio | mor (g) | $\mathrm{m}_{\mathrm{r}}(\mathrm{g})$ | $\mathbf{m o p}_{\text {op }}(\mathrm{g})$ | $\mathrm{m}_{\mathrm{p}}(\mathrm{g})$ | $\begin{aligned} & \rho \mathrm{r} \\ & \left(\mathrm{~g} / \mathrm{cm}^{3}\right) \end{aligned}$ | $\rho_{\mathrm{p}}\left(\mathrm{~g} / \mathrm{cm}^{3}\right.$ ) |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 0:100 |  |  |  |  |  |  |
| 20:80 |  |  |  |  |  |  |
| 40:60 |  |  |  |  |  |  |
| 60:40 |  |  |  |  |  |  |
| 80:20 |  |  |  |  |  |  |
| 100:0 |  |  |  |  |  |  |

Since in columns II and VI of the Table there are quantities referring to distilled water, in these columns a single value is entered, common to the six compositions.
" $\rho_{p}$ " is graphically represented as a function of the volumetric percentage composition of the ethanolwater mixture, according to the model in Figure 10. The broken line represents the theoretical plot of the density vs. composition for the case that volume contraction would not occur when preparing the ethanol-water mixture. The graph resulting from the experimental data is located above this line, a sign that when ethanol is mixed with water, the density has higher values than theoretically predicted, so volume contraction occurs.


Fig. 4 Model for graphical representation of results

## DETERMINATION OF LIQUID DENSITY USING DENSIMETERS

## INTRODUCTION

## DENSITY

It is a physical quantity that characterizes the mass of the volume unit. It is generally noted with the Greek letter $\rho$ (rhô), as it is also used by the International Bureau of Weights and Measures. Density is by definition the ratio between mass and volume $\rho=\frac{m}{V}$

Where:
m - the mass of the studied material
V- the volume of the studied material
Density is synonymous with absolute density and specific density

The unit of measurement in the International System is: kilogram per cubic meter: (kg $\cdot \mathrm{m}^{-}$ ${ }^{3}$ or $\left.\mathrm{kg} / \mathrm{m}^{3}\right):[\rho]_{S I}=\frac{k g}{m^{3}}$ The unit of measurement in the CGS System is: gram per cubic centimeter: $\left(\mathrm{g} \cdot \mathrm{cm}^{-3}\right.$ or $\left.\mathrm{g} / \mathrm{cm}^{3}\right):[\rho]_{C G S}=\frac{g}{\mathrm{~cm}^{3}}$

Water is often used as a reference liquid for calculating the density of liquids and solids.
The density of water is $1000 \mathrm{~kg} / \mathrm{m}^{3}$ (or $1 \mathrm{~kg} / \mathrm{dm}^{3}$ or $1 \mathrm{~kg} / \mathrm{l}$ or $1 \mathrm{~g} / \mathrm{cm}^{3}$ ). This is the density of water at $3.98^{\circ} \mathrm{C}$ and normal atmospheric pressure. Indeed, water has a peculiarity among liquids: its density is maximum, not at the melting temperature $\left(0^{\circ} \mathrm{C}\right)$ like most liquids, but at $3.98{ }^{\circ} \mathrm{C}$. In metrology, taking as a reference a physical property that passes through an extreme it is very interesting what happens in the vicinity of this extreme: the physical property varies very little. Thus, in the vicinity of $3.98^{\circ} \mathrm{C}$, an inaccuracy in determining the exact temperature affects very little the density of water and therefore the result of the density measurement
Relative density: $d=\frac{\rho_{\text {solid }}}{\rho_{\text {water }}} \quad$ or $\quad d=\frac{\rho_{\text {liquid }}}{\rho_{\text {water }}}$
Where $\rho_{\text {material }}$ is the density of the considered material (liquid or solid) and $\rho_{\text {water }}$ is the density of water ( $1000 \mathrm{~kg} / \mathrm{m}^{3}$ ).

To correct the values measured at temperatures other than the reference temperature, it is necessary to use the curve of variation of the density of water as a function of temperature, in order to deduce the density of the materials at the temperature of the measurements.

| Liquids | $\boldsymbol{\rho ( \mathbf { k g } / \mathbf { m } ^ { 3 } )}$ |
| :--- | :--- |
| Acetone | 790 |
| Acetic acid | 1049 |
| liquid nitrogen at $-195^{\circ} \mathrm{C}$ | 810 |
| bromine at $0^{\circ} \mathrm{C}$ | 3087 |
| Water to $4^{\circ} \mathrm{C}$ | 1000.0 |
| Sea water | 1000 to 1032 |
| benzine | 750 |
| ethanol | 789 |
| ether | 710 |
| diesel | 1260 |
| glycerin | 150 |
| liquid helium at $-269^{\circ} \mathrm{C}$ | 920 |
| olive oil | 70 |
| liquid hydrogen at $-252{ }^{\circ} \mathrm{C}$ | 70 |
| liquid oxygen at $-184^{\circ} \mathrm{C}$ | 1140 |
| milk | 1030 |

## ARCHIMEDES' THEOREM

A body immersed in a liquid is pushed from the bottom up with a force equal to the weight of the volume of fluid displaced by the body.

The theorem applies to bodies immersed in liquids at rest.

The theorem applies if the immersion liquid and the immersed body are at rest. It is also necessary to replace the immersed body with immersion liquid without disturbing the balance. A counter-example is the cap of a bathroom filled with water: if it is replaced by water, it is clear that the bathroom is empty and that the fluid is no longer at rest. The theorem does not apply because we are in a case where the plug is not completely wetted by the liquid and does not pass its free surface. In a uniform gravitational field, the Archimedean force, noted with the symbol $\mathrm{F}_{\mathrm{A}}$, is given by the formula:

$$
\vec{F}_{A}=-M_{f} \vec{g}
$$

Where $\mathrm{M}_{\mathrm{f}}$ is the mass of the fluid contained in the displaced volume V and g is the gravitational acceleration.

If the mass of the displaced fluid is expressed as a function of density, the Archimedean force has the form: $\vec{F}_{A}=-\rho \cdot V \cdot \vec{g}$

In scalar form it is written as follows: $\left\|\vec{F}_{A}\right\|=-\rho \cdot V \cdot g$


Figure 1: Schematic representation of Archimedes' law

The Archimedean force $\mathrm{F}_{\mathrm{A}}$ and the weight G can be in one of the 3 cases presented (from left to right):
Case 1. $\mathbf{G}>\mathbf{F}_{\mathbf{A}}$ (the weight $G$ is greater than the Archimedean force $\mathrm{F}_{\mathrm{A}}$ )
In this case the apparent weight appears: $\mathrm{G}_{\mathrm{a}}=\mathrm{G}-\mathrm{F}_{\mathrm{A}}$
The body falls to the base of the vessel
Case 2. $\mathbf{G}=\mathbf{F}_{\mathbf{A}}$ the body is in equilibrium inside the liquid.
Case 3. $\mathbf{G}<\mathbf{F}_{\mathbf{A}}$ (The Archimedean force $\mathrm{F}_{\mathrm{A}}$ is greater than the weight $G$ )
In this case the ascension force appears: $\mathrm{Fa}=\mathrm{F}_{\mathrm{A}}-\mathrm{G}$
$F_{\mathrm{p}}=\rho \boldsymbol{V} g$
$F_{\mathrm{A}}=\rho_{\mathrm{f}} \boldsymbol{V} \boldsymbol{g}$
$\boldsymbol{G} / \boldsymbol{F}_{\mathrm{A}}=\boldsymbol{\rho}_{\mathrm{s}} / \boldsymbol{\rho f}_{\mathrm{f}}$
The force ratio is equal to that of the densities

- If the density of the solid material is greater than that of the fluid $G>F_{A}$, the solid falls to the bottom of the vessel
- If the density of the solid material is lower than that of the fluid $G<F_{A}$, the solid floats on the surface of the liquid from the vessel
- If the density of the solid material is equal to that of the fluid $G=F_{A}$, the solid remains in equilibrium inside the liquid from the vessel.


## MEASURING THE DENSITY OF LIQUIDS USING THE DENSIMETER

The densimeter consists of a glass cylinder filled with air, which has weights at the base and is graduated. It enters into the liquid to be measured more or less, depending on its density.

The density of the liquid in which it is immersed is read directly on the graduation that is in contact with the free surface of the liquid.

The hydrometer is mainly used to measure the density of grape must, barley, sugar syrup or to evaluate the alcohol content. The hydrometer is also used in sedimentometry laboratory tests to calculate the distribution of microscopic particles in a soil sample.

It is also used in the ceramic industry as a means of controlling the density of slips, mainly for the density of glazes or slips that will cover tiles or bricks. This type of measurement has the advantage of being easy to implement and allowing very fast control, but is less accurate than density measurement if a precise volume is weighed using standard 100 ml containers


The Fahrenheit hydrometer looks like a glass tube filled with air in the middle, with mercury at the bottom and contains a container at the top. When the hydrometer is immersed in a liquid of unknown density, it floats. The marked masses are then added to the upper container until the center of the hydrometer is completely immersed. The total mass added is the greater Argimedic force imposed on the hydrometer by the liquid,
i.e. the latter's density is greater.

Figure 2. Variants of densimeters, hydrometers, alcohol meters

The alcohol meter is used to measure the degree of alcoholic strength starting from the density of the liquid was defined by Louis Joseph Gay-Lussac in 1824, for the purpose of taxation on alcohol. A hydrometer directly graduated in degrees of alcohol is immersed in the liquid. The measured value is corrected according to the temperature of the liquid. The density tables and their temperature corrections were established experimentally by Gay-Lussac. This technique is accurate only if the liquid contains only water and alcohol in solution, hence the need for previous distillation, if rigorous data is required for commercial or fiscal needs.

Small hydrometers graduated in degrees, alcohol meters or alcohol scales are used by professionals and controllers to roughly and quickly measure the alcoholic strength of a product.

Today, the Malligand Ebulliometer, which appeared in 1875, is used to measure wine strength, which is a miniature still equipped with a precision thermometer. The boiling point of an ethanolwater solution depends on their ratios. The temperature read is converted to alcohol by reading a table. The accuracy of the measurement depends on the concentration of other substances, such as sugars and carbonic acid, which change the boiling point. In the wine industry, the expression "Malligand degree" has long been used for "\% vol. at $15^{\circ} \mathrm{C}$ ".

Appeared in 1967, an electronic densimeter allows the measurement of TAV (degree of alcohol) from an aqueous alcoholic solution, which requires a distillation before measurement, if the samples have a complex matrix (wine, cognac, Armagnac, liqueurs .. .).

Oscillating tube density measurement (electronic density meter) became the reference standard method in the 1990s for many producers and specialists in wines, spirits and more recently for ethanol as well as precision fuel (TAV $+/-0.01 \% \mathrm{v} / \mathrm{v}$ ) and ease of use.

## DETERMINATION OF THE DENSITY OF A LIQUID USING A DENSIMETRE

Determination of the density of a liquid with a hydrometer is based on Archimedes ' law regarding the upward force exerted on a body immersed in a liquid. Figure 3 represents a densimeter. This is an instrument made of glass in which three portions are distinguished: the weight " G ", the volume body "V" and the rod/bar provided with divisions "T" (Figure 3-A). The studied liquid is introduced into a glass and afterwards, the densimeter is introduced into the liquid. The hydrometer floats in the liquid and the depth to which it sinks depends on the density of the liquid. The depth is read on the divisions made on the " T " bar, which correspond to the density values of the liquid (Figure 3-B).

When performing the determination, it is important to control the temperature of the liquid. For this reason, a thermometer is also included in some densimeters.


To cover a sufficiently wide range of density values, sets of densimeters are used, each calibrated for a specific density range. While performing the determination of the density of a concrete liquid, choose the densimeter from the set which, inserted in the liquid, floats without touching the bottom of the vessel.

In practice, the density of the following liquids is determined: ethanol $\left(\mathrm{C}_{2} \mathrm{H}_{5} \mathrm{OH}\right)$, carbon tetrachloride $\left(\mathrm{CCl}_{4}\right)$, acetone $\left(\left(\mathrm{CH}_{3}\right)_{2} \mathrm{CO}\right)$. When performing the experiment, it is important that the densimeter does not touch the wall of the vessel in which the studied liquid is placed.

Because the densimeters are calibrated by the manufacturer, determination of the densities is not dependent on any reference liquid.

Figure 3. Simple densimeter

## PROCEDURE FOR THE PRACTICAL WORK - STEPS

1. 3 measurements are made for each liquid ( NaCl and KCl solutions) with concentrations of $\mathbf{0 . 9 \%}, 1 \%, 5 \%, 6 \%$ and $8 \%$
2. Average values, errors and actual values will be calculated
3. Calculations are made for the relative density - $d$
4. Average values, errors and actual values will be calculated
5. The obtained values are entered in the table
6. The graphic representation will be made, respecting the function: $\rho=\rho(\mathrm{c})$

Table: densities of $\mathbf{N a C l}$ solutions

| No | Sol. used | $\begin{aligned} & \rho \\ & (\mathrm{g} / \mathrm{l}) \end{aligned}$ | $\bar{\rho}$ (g/l) | $\begin{aligned} & \mu_{\rho} \\ & (\mathrm{g} / \mathrm{l}) \end{aligned}$ | $\rho_{\text {real }}(\mathrm{g} / \mathrm{l})$ | $d=\frac{\rho_{\text {liquid }}}{\rho_{\text {water }}}$ | $\bar{d}$ | $\mu_{\text {d }}$ | $\mathrm{d}_{\text {real }}$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 1 | $\mathrm{H}_{2} \mathrm{O}$ |  |  |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |  |  |
| 1 | $\begin{aligned} & \mathrm{NaCl} \\ & 0.9 \% \end{aligned}$ |  |  |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |  |  |
| 1 | $\mathrm{NaCl} 1 \%$ |  |  |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |  |  |
| 1 | $\mathrm{NaCl} 5 \%$ |  |  |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |  |  |
| 1 | NaCl 6\% |  |  |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |  |  |
| 1 | $\mathrm{NaCl} 8 \%$ |  |  |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |  |  |

Table: densities of $\mathbf{K C l}$ solutions

| $\begin{aligned} & \text { No.c } \\ & \text { rt } \end{aligned}$ | Sol. used | $\rho$ <br> (g/l) | $\bar{\rho}$ (g/l) | $\mu_{\rho}$ (g/l) | $\rho_{\text {real }}$ <br> (g/l) | $d=\frac{\rho_{\text {liquid }}}{\rho_{\text {water }}}$ | $\bar{d}$ | $\mu_{\text {d }}$ | $\mathrm{d}_{\text {real }}$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 1 | $\mathrm{H}_{2} \mathrm{O}$ |  |  |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |  |  |
| 1 | $\begin{aligned} & \mathbf{K C l} \\ & 0.9 \% \end{aligned}$ |  |  |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |  |  |
| 1 | KCI 1\% |  |  |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |  |  |
| 1 | KCI 5\% |  |  |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |  |  |
| 1 | KCl 6\% |  |  |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |  |  |
| 1 | KCI 8\% |  |  |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |  |  |

## DETERMINATION OF LIQUID DENSITY USING THE ABBE METHOD

## THEORETICAL INTRODUCTION

In wave physics，refraction means the change in the wave propagation direction when it crosses the separation surface between two media（optical，acoustic or even seismic）

Refraction means a change in orientation：
－Wave front－is the line that describes a light wave（in physical optics or seismology）
－Of the beam：the direction of wave propagation is perpendicular to that of the wave front The two changes of orientation are equivalent in the case of refraction，however，we prefer the first to explain the phenomenon，and the second to quantify the phenomenon

Light is deflected when it passes from one transparent medium to another（for example，from air to water，or vice versa）．This phenomenon is noticed when we look to a straw in a glass and it appears broken．This＂division＂is possible due to＂refraction＂．

Light is called＂refracted＂，and the property that characterizes different transparent media is ＂refraction＂．

Diffuse refraction can be distinguished from perfect refraction．
－in the case of diffuse refraction，the incident ray is separated into a multitude of rays in the second medium of propagation；
－in the case of perfect refraction，only one ray appears，in the second medium．


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Diffuse refraction generally occurs in nature．For perfect refraction，the optical surface（the separation surface between the circles）must be perfectly smooth，and the second medium must be perfectly transparent and amorphous or monocrystalline．

The first concepts discovered in experimental optics were those related to the phenomena of reflection and refraction．

Refraction was already well known by Ptolemy

The first to mention the law of refraction was Ibn Sahl．（c．940－1000）

## Snell-Descartes law



In optics, each transparent medium is characterized by the refractive index ni Diopter is the surface of separation between two media.

The laws of refraction (the second law), established by Snell and Descartes, can explain the quantitative phenomenon, states that:

- The incident, normal and refracted rays are coplanar.
- The refractive index noted with $\mathbf{n}$ is the ratio between the sin of the angle of incidence $\mathbf{i}$ and the $\sin$ of the angle of refraction $\mathbf{r}$
$\mathbf{n}_{\mathbf{1}} \cdot \boldsymbol{\operatorname { s i n }} \mathbf{i}=\mathbf{n}_{\mathbf{2}} \cdot \boldsymbol{\operatorname { s i n }} \mathbf{r}$, if $\mathrm{n}_{1}=1$ and $\mathrm{n}_{2}=\mathrm{n}$ then:
- $n=\frac{\sin i}{\sin r}$

Then it is observed that:

- The higher the refractive index $n_{2}$, the more the refracted ray approaches the normal, and vice versa;
- If the refractive index $\mathrm{n}_{2}$ is lower than $\mathrm{n}_{1}$ (for example, the ray passing from glass to air), the socalled "critical angle" can be exceeded, beyond which a total reflection is observed.


## ABBE REFRACTOMETER DESCRIPTION:

The refractometer is a device used to determine the unknown refractive index of some solutions. It has been used since the 19th century.

In principle, all solutions now have a unique refractive index that can be used for its identification.

The refractive index of organic compounds is in the range $\mathrm{n} \in(1.08,1.25)$
The Abbe refractometer is successfully used to determine the refractive index of organic compounds, solutions, food products, pharmaceuticals and serum protein concentration.

The refractive index is measured by aligning the reticle in the telescope with the line corresponding to the total reflection (see the demarcation line between the bright and dark areas of the telescope).

This line is moved by rotating the prism. The reading is carried out at a constant temperature, which is important, so the prisms are often enclosed in a water jacket, which can be connected to a thermostatic bath.

The Carl Zeiss company and the house of Abbe originated the refractometer for much of its history. Ernst Abbe published his book in 1874: »Neue lui APPARATE ....«, in which he discussed the described theory and the instruments that he had developed for measuring the refractive index using prisms and total reflection. He first described the Abbe refractometer invented in 1863 to determine the refractive index of fluids with a refractive index in the range ( $\mathrm{n}=1.03-1.07$ ). This instrument includes the original Amici prisms and is essentially the same as the modern Abbe refractometer, but without the temperature gain. Initially, Carl Zeiss produced this instrument for internal use and for private customers. Only after 1884 Abbe offered in catalogs for Zeiss microscopes but without illustrations. Carl Zeiss published its first instrument catalog in 1893, following the formation of a new optical instrument department under Dr. C. Pulfrich. The Abbe refractometer was the prominent figure in this catalog, both with and without water jacket prisms.

where:

1. plate for lighting - where the incident light falls
2. calibration screw
3. focus adjustment
4. ocular
5. contrary calibration screw

## PROCEDURE FOR THE PRACTICAL WORK - STEPS

1. Place the front end of the refractometer in the direction of the light
2. Zero adjustment: Open the refractometer plate towards the bright region
3. Place 2-3 drops of distilled water on the main prism.
4. Close the lid of the prism - close the access to the light and press gently, so that the water spreads over the entire surface of the prism without air bubbles or dry spots.
5. Allow the sample to remain on the prism for approximately 30 seconds.
6. Then the calibration of the refractometer is carried out so that the boundary light coincides with the zero line corresponding to darkness.
7. The setting of the refractometer with the temperature compensation function to be done below the values of $20^{\circ} \mathrm{C}$ which is the room temperature. When the temperature of the room or working environment (without sample) changes by more than $5^{\circ} \mathrm{C}$, it is recommended to mention the calibration accuracy.
8. The work procedure is done after calibration and is basically done in the same way as calibration.
9. The refractometer plate is opened.
10. Clean the surface of the prism with a soft cotton cloth.
11. Place 2-3 drops of solution on the main prism.
12. Close the daylight plate and press lightly, then read the value on the appropriate dark limit scale.
13. The reading is the measured value of the density of the solution.
14. After the measurement, immediately clean the test sample from the surface of the prism and the cover plate with a damp cotton cloth. After drying, it must be closed perfectly.

## DETERMINATION OF THE CONCENTRATION OF THE SOLUTION

1. Prepare different solutions of propylene glycol $\mathrm{C}_{3} \mathrm{H}_{8} \mathrm{O}_{2}$ with distilled water with concentrations between 10-80\%
2. Read the density value on the device (Abbe refractometer) on the scale on the left
3. Using the formula, $c=\frac{m(g)}{M . V(l)}=\frac{\rho}{M}$ the concentration can be determined. Where M is the molar mass of propylene glycol $\mathrm{C}_{3} \mathrm{H}_{8} \mathrm{O}_{2}$ that is: $\mathrm{M}_{\mathrm{C} 3 \mathrm{H} 8 \mathrm{O} 2}=3 \cdot 12+8 \cdot 1+2 \cdot 16=76 \mathrm{~g} / \mathrm{mol}$, thus MC3H8O2 $^{\mathbf{~}} \mathbf{7 6} \mathbf{~ g} / \mathrm{mol}$
4. Enter the data in the table below
5. The graphic representation will be made, respecting the function: $\rho=\rho(\mathrm{c})$
6. Enter the data in the table below
7. Read the values for the freezing temperature on the middle scale and write the corresponding data in the table

Table 1: variation of density with concentration and freezing temperature for ethylene glycol solutions

| Crt. <br> no. | $\mathrm{c}_{\mathrm{p}}(\%)$ | $\mathrm{c}(\mathrm{M})$ | $\rho\left(\frac{\mathrm{kg}}{\mathrm{l}}\right)$ | $\mathrm{t}^{\mathrm{o}} \mathrm{C}$ |
| :--- | :--- | :--- | :--- | :--- |
| 1 | 10 |  |  |  |
| 2 | 20 |  |  |  |
| 3 | 30 |  |  |  |
| 4 | 40 |  |  |  |
| 5 | 50 |  |  |  |
| 6 | 60 |  |  |  |
| 7 | 70 |  |  |  |
| 8 | 80 |  |  |  |

$$
c_{M}=\frac{\rho \cdot c_{p} \cdot 10}{M}
$$

The experiment can be performed with different solutions of clarithromycin, vitamin C solutions, etc.

## DETERMINATION OF SURFACE TENSION OF LIQUIDS

## THEORETICAL INTRODUCTION

In the case of a liquid in contact with a gaseous medium, a molecule inside it is uniformly stressed in all directions by the other molecules, so that the resultant of these forces is zero, on the other hand, a molecule located at the liquid-gas separation surface will be subjected to an attraction force towards the inside of the liquid, because the gas molecules being less numerous than the liquid ones, the attraction forces exerted by them will be smaller. These molecules constitute the superficial layer of the liquid and the force with which the surface unit is attracted from the superficial layer towards the inside of the liquid is called molecular pressure. So the superficial layer has a higher potential energy than the molecules inside, being proportional to the surface of the upper layer. According to the minimum condition for this surface energy, the surface of a liquid always tends to the shape that gives it a minimum area. So the surface of a liquid behaves like a stretched membrane, and the force that tends to measure the surface is called the surface tension force. If we noted with dL the mechanical work done to increase the surface of the membrane by dS , we write:

$$
\begin{equation*}
\mathrm{dL}=-\sigma \mathrm{dS} \tag{1}
\end{equation*}
$$

where $\sigma$ is the SURFACE TENSION COEFFICIENT.
If the process carried out with the film is a reversible and isothermal process (it takes place at the ambient temperature), this mechanical thing represents precisely the free energy of the membrane:

$$
\begin{align*}
& -\mathrm{dE}=\mathrm{d} \mathrm{~L}  \tag{2}\\
& \mathrm{dE}=\sigma \mathrm{dS} \tag{3}
\end{align*}
$$

Integrating relation (3) between $\mathrm{S}=0$ and $\mathrm{S}=\mathrm{S}$ we obtain:

$$
\begin{equation*}
E=\sigma S \tag{4}
\end{equation*}
$$

Consequently, the surface tension force will be of the form:

$$
\begin{equation*}
\mathrm{F}=\sigma \cdot \mathrm{I} \tag{5}
\end{equation*}
$$

where 1 represents the length of the contour on which the surface tension force F is applied.
From relation (4) we can deduce that the surface tension coefficient represents the free energy related to the area loss from the surface of the superficial membrane.

Consider a capillary tube open at one end. The detachment of a drop that originates at the end of a capillary tube occurs when its weight overcomes the surface tension force.

In this case we can write:
or

$$
\begin{align*}
& \mathrm{m} \cdot \mathrm{~g}=2 \cdot \pi \cdot \sigma  \tag{6}\\
& \mathrm{~V}_{1} \cdot \rho \cdot \mathrm{~g}=2 \cdot \pi \cdot \mathrm{r} \cdot \sigma
\end{align*}
$$

where $\mathrm{V}_{1}$ is the volume of a drop, r is the radius of the capillary tube, so the surface tension coefficient will be as follows:

$$
\begin{equation*}
\sigma=\frac{V_{1} \rho g}{2 \pi r} \tag{8}
\end{equation*}
$$

The volume $\mathrm{V}_{1}$ of a drop can be determined by knowing the volume V of the liquid and the number n of drops that arise from this volume:

$$
\begin{equation*}
V_{1}=\frac{V}{n} \tag{9}
\end{equation*}
$$

If we have a known liquid characterized by the surface tension coefficient $\sigma_{o}$ and the tension $\rho_{o}$, we write the relation (10). By dividing relation (8) by (10) we can write relation (11).

$$
\begin{align*}
\sigma_{o} & =\frac{V \rho_{o} g}{2 \pi r n_{o}}  \tag{10}\\
\sigma & =\sigma_{o} \frac{\rho n_{o}}{\rho_{o} n} \tag{11}
\end{align*}
$$

The drop method gives us relative measurements. It represents the advantage that with this method, the measurements do not depend on the limit angle of the liquid in relation to the capillary material and the method can also be used for viscous liquids (liquids with higher viscosity). Also these measurements are as simple and precise.

In this paper, we propose to determine the coefficient of surface tension for different liquids, with the help of the stalagmometer.

## SODIUM DODECYL SULFATE - SDS

Sodium dodecyl sulfate (SDS or NaDS), sodium lauryl sulfate or sodium lauryl (SLS) is an organic compound with the formula $\mathrm{CH}_{3}\left(\mathrm{CH}_{2}\right)_{11} \mathrm{OSO}_{3} \mathrm{Na}$.

This is an anionic surfactant used in many cleaning and hygiene products.
The salt is composed of an inorganic sulfate with a chain of 12 carbon atoms attached to a sulfate group, giving the material the properties required for amphiphilic detergents. Derived from coconut and palm oils, it is a common element in many household cleaners.

PRODUCTION: SDS is synthesized by treating lauryl alcohol with sulfur trioxide gas, oils, where a hydrogen atom can produce lauryl sulfate. The method frequently practiced industrially uses sulfur trioxide as gas. The product obtained is then neutralized by adding sodium hydroxide or sodium carbonate. Lauric alcohol is in turn generally derived from hydrolyzed coconut oil and palm oil, releasing the fatty acids, followed by hydrogenation. Because of this method of synthesis, commercial samples of SDS are often a mixture of other alkyl sulfates, with dodecyl sulfate as the main component. SDS is commercially available as powders and granules. It appears that the degranule form dissolves faster than the powder in water.

APPLICATIONS: SDS is primarily used as laundry detergents with many cleaning applications. SDS is an extremely effective surface active agent and is used in a task that requires the removal of stains and oily residues.

For example, it is found in higher concentrations in industrial products, including engine degreasers, floor cleaners and car wash soaps. It is found in toothpastes, shampoos, shaving creams, bath foam solutions.

It can be used to aid in cell lysis during DNA extraction and protein breakdown in SDS-PAGE. Sodium lauryl sulfate, scientifically called sodium dodecyl sulfate (SDS) or Duponol, is frequently used in the preparation of proteins for electrophoresis according to the SDS-PAGE technique. This compound works by disrupting the non-covalent bonds of proteins, denaturing and causing the molecules to lose their native shape (conformation).

This new negative charge is substantially greater than the original charge of this protein. The electrostatic repulsion that is created by binding to SDS causes the proteins to unfold into a rod-like shape, eliminating differences in shape. Electrostatic repulsion acts as a separation factor in the gel.

Sodium lauryl sulfate is probably the most studied anionic surfactant. Like all surfactants in detergents, sodium lauryl sulfate strips skin of oils and can irritate skin and eyes.

The critical micelle concentration (CMC) in pure water at $25{ }^{\circ} \mathrm{C}$ is 8.2 mM , and the aggregation number at this concentration is generally considered to be about 62.

The micellar ionization fraction (of $\alpha$ ) is approximately 0.3 (or 30\%).
Aqueous solutions of SDS are also popular for dispersing or suspending nanotubes such as carbon nanotubes (CNTs). However, sodium dodecylbenzene sulfonate (NaDDBS) and Triton X-100 are reported to be better dispersed on carbon walls.

Sodium lauryl sulfate is also used in hemoglobin analysis. The hydrophobic group of SLS acts on the globin subunit, causing a conformational change. Biological tissues treated with sodium dodecyl sulfate are subjected to quasi-transparent optical microscopy that allows to see through them.

SDS is a potentially effective topical bactericide that can also inhibit and possibly prevent infection with various enveloped and non-encapsulated viruses, such as herpes simplex virus, HIV, and Semliki Forest virus.

Pesticides: Some species of flies, ants and moths are attracted to sodium lauryl sulfate, even in small amounts, and then killed. This is the main active ingredient in most store-bought fruit fly repellants. fruit flies are often attracted to it.

Taste disturbances: Sodium lauryl sulfate temporarily reduces the perception of sweetness, an effect commonly seen after recent use of toothpastes containing this ingredient.

Toxicology studies have shown that SDS is not carcinogenic when applied directly to the skin or consumed. A review of the scientific literature said "SLS [SDS] was negative in an Ames Bacterial Mutation Test, a gene mutation and sister chromatid exchange test in mammalian cells, and an in vivo test on the micronuclei of mice. The negative results of the in vitro and in vivo studies indicate that SDS does not interact with DNA. The only study conducted on the carcinogenicity study said "SDS was not carcinogenic in the Beagle". In the case studied, details were lacking and the data had limited significance "(Agner T (1991)." Susceptibility to atopic dermatitis of patients with irritant dermatitis is caused by sodium lauryl sulfate "Acta Derm Venereol 71 (4): ... 296-300)
Sensitivity : It has been shown to irritate the skin of the face during prolonged and constant exposure (over one hour) in young adults. SDS can worsen skin problems in people with chronic skin hypersensitivity. Some people are more affected than others. In animal studies SDS appears to irritate the skin and eyes.

Canker sores: There have been several studies on how SDS toothpaste affects the recurrence of canker sores, known in some countries as white sores. The results of these studies were in conflict. In 1994, a preliminary central study showed that patients experienced a significantly higher number
of aphthous ulcers after using toothpaste containing SLS than toothpaste without SLS. However, a study published in 1999 failed to find statistically significant differences. A 2012 study found no significant difference in the number of ulcers, but did not find a significant difference in the duration and scores of ulcer-related pain. Only the 2012 Study of patients using a toothpaste without SLS showed faster ulcer healing and less ulcer-related pain, on average, than patients using toothpaste containing SLS.

Interaction with fluoride: Some studies have suggested that SLS in toothpaste may decrease the effectiveness of fluoride in preventing tooth decay (cavities). This may be due to the interaction with SLS deposition of fluoride on tooth enamel.

## EXPERIMENTAL ASSEMBLY

THE STALAGMOMETER or BURETTE (figure 1) consists of a reservoir of a determined volume, delimited by two landmarks, finished with a capillary tube, through which the liquid flows, drop by drop.


Figure 1. Schematic representation of stalagmometer

Water is introduced into the stalagmometer. Count the water droplets formed while the water flows between the two landmarks (in a given volume). The operations are repeated for the other substances,
which are being studied. 5 determinations are made for each substance, whose density is known and whose surface tension coefficient is determined with relation (11). The results are shown in table 1.

## PROCEDURE FOR THE PRACTICAL WORK - STEPS

1. Water is introduced into the stalogmometer
2. The drops that forms from a given volume (between 2 marks) of liquid are counted
3. The operation is repeated for all other substances used
4. 3 measurements are made for each substance
5. The surface tension coefficient is determined with formula (11)
6. The data are entered in table 1 and 2 , respectively.

## EXPERIMENTAL DATA PROCESSING

Knowing the density of the solutions, the surface tension coefficient of the fluids is determined with the relation: $\sigma=\sigma_{0} \frac{\rho n_{0}}{\rho_{0} n}$

The results will be introduced in table 1.
For each substance, the average value is calculated - $\bar{\sigma}$
Also the error will be calculated with the following formula $\mu=\sqrt{\frac{\sum\left(\sigma_{i}-\bar{\sigma}\right)}{n(n-1)}}$
and the interval of real values $\sigma_{\text {real }} \in(\bar{\sigma}-\mu, \bar{\sigma}+\mu) 10^{-3} \frac{\mathrm{~N}}{\mathrm{~m}}$

Table 1. Variation of surface tension with density

| No. <br> crt. | Substance | $\rho\left(\mathrm{kg} / \mathrm{m}^{3}\right)$ | $n$ <br> $($ no. drops $)$ |  | $\sigma\left(10^{-3} \mathrm{~N} / \mathrm{m}\right)$ | $\sigma$ |
| :--- | :--- | :--- | :--- | :--- | :--- | :--- |
| 1 |  |  |  |  |  |  |

Table 2: Variation of surface tension with concentration

| No crt | Subst. | $\begin{aligned} & \rho \\ & \left(\mathrm{kg} / \mathrm{m}^{3}\right) \end{aligned}$ | $n$ (no. drops) | $\begin{aligned} & \sigma \\ & \left(10^{-3} N / m\right) \end{aligned}$ | $\bar{\sigma}\left(10^{-3} \mathrm{~N} / \mathrm{m}\right)$ | $\begin{gathered} \mu \\ \left(10^{-3} N / m\right) \end{gathered}$ | $\sigma_{\text {real }}$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 1 | water$\left(\mathrm{H}_{2} \mathrm{O}\right)$ | 1000 |  | $\sigma_{0}=72.8$ | - | - |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |
| 1 | $\begin{aligned} & \text { SDS } \\ & 0.5 \% \end{aligned}$ |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |
| 1 | $\begin{aligned} & \text { SDS } \\ & 1 \% \end{aligned}$ |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |
| 1 | $\begin{aligned} & \text { SDS } \\ & 2 \% \end{aligned}$ |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |
| 1 | $\begin{aligned} & \text { SDS } \\ & 5 \% \end{aligned}$ |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |
| 1 | $\begin{aligned} & \text { SDS } \\ & 6 \% \end{aligned}$ |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |
| 1 | $\begin{aligned} & \hline \text { SDS } \\ & 9.7 \% \end{aligned}$ |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |

## CONCLUSIONS

Conclusions are drawn regarding the variation of the surface tension with the density, respectively with the concentration of the studied substance.

The interval of real values for the studied liquids is calculated according to the formula:
$\sigma_{\text {real }} \in(\bar{\sigma}-\mu, \bar{\sigma}+\mu) 10^{-3} \frac{\mathrm{~N}}{\mathrm{~m}}$

Table 3: Variation of surface tension with concentration:
MICELLAR SOLUTION

| $\begin{aligned} & \mathrm{No} . \\ & \mathrm{crt} \end{aligned}$ | C (\%) | $\begin{aligned} & \hline \rho \\ & \left(\mathrm{kg} / \mathrm{m}^{3}\right) \end{aligned}$ | $n$ <br> (no. drops) | $\sigma\left(10^{-3} \mathrm{~N} / \mathrm{m}\right)$ | $\sigma\left(10^{-3} \mathrm{~N} / \mathrm{m}\right)$ | $\begin{aligned} & \mu \\ & \left(10^{-3} \mathrm{~N} / \mathrm{m}\right) \end{aligned}$ | $\sigma_{\text {real }}$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 1 | water$\left(\mathrm{H}_{2} \mathrm{O}\right)$ | 1000 |  | $\sigma_{0}=72.8$ | - | - |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |
| 1 | 10 | 988.4 |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |
| 1 | 20 | 989.2 |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |
| 1 | 30 | 990 |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |
| 1 | 40 | 990.4 |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |

## CONCLUSIONS

Conclusions are drawn regarding the variation of the surface tension with the density, respectively with the concentration of the studied substance.

The interval of real values for the studied liquids is calculated according to the formula: $\sigma_{\text {real }} \in(\bar{\sigma}-\mu, \bar{\sigma}+\mu) 10^{-3} \frac{\mathrm{~N}}{\mathrm{~m}}$

Table 4: Variation of surface tension with concentration: GLYCERIN

| No. <br> crt | C(\%) | $\rho\left(\mathrm{kg} / \mathrm{m}^{3}\right)$ | $n$ <br> (no drops) | $\sigma\left(10^{-3} \mathrm{~N} / \mathrm{m}\right)$ | $\bar{\sigma}\left(10^{-3} \mathrm{~N} / \mathrm{m}\right)$ | $\begin{gathered} \mu \\ \left(10^{-3} N / m\right) \end{gathered}$ | $\sigma_{\text {real }}$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 1 | $\begin{aligned} & \text { water } \\ & \left(\mathrm{H}_{2} \mathrm{O}\right) \end{aligned}$ | 1000 |  | $\sigma_{0}=72.8$ | - | - |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |
| 1 | 10 |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |
| 1 | 20 |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |
| 1 | 30 |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |
| 1 | 40 |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |
| 1 | 50 |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |

## CONCLUSIONS

Conclusions are drawn regarding the variation of the surface tension with the density, respectively with the concentration of the studied substance.

The interval of real values for the studied liquids is calculated according to the formula: $\sigma_{\text {real }} \in(\bar{\sigma}-\mu, \bar{\sigma}+\mu) 10^{-3} \frac{\mathrm{~N}}{\mathrm{~m}}$

# MEASUREMENT OF DYNAMIC AND KINEMATIC VISCOSITY USING THE HOPPLER <br> <br> VISCOSIMETER 

 <br> <br> VISCOSIMETER}

## THEORETICAL INTRODUCTION

VISCOSITY
At not too high speeds, the fluid flow is linear (in parallel layers), thus, the current lines are well determined and do not intersect anywhere, each fluid particle always remains inside the same current tube. At high speeds, the movement of the fluid becomes TURBULENT, irregular, the portions of the fluid mix (the flow of the ideal fluid with vortices). If the fluid layers slide against each other, internal friction or viscosity forces appear between them. The layer with the lower speed will brake the layer with the higher speed and vice versa the layer with the higher speed will increase the speed of the layer with the lower speed.

The appearance of these forces, located in the working planes, is explained by the variation of the momentum of the layers due to the passage of molecules from one layer to another. We will assume that the fluid flow direction is the same and that the flow speed varies as a mode only in the direction perpendicular (transversal) to the flow direction.

The example shows that the internal friction force that appears in the sliding plane per surface unit is proportional to the degree of speed (Newton's law):

$$
\begin{equation*}
\mathfrak{I}=\frac{d F}{d S}=\eta \frac{d v}{d z} \tag{1}
\end{equation*}
$$

thus $\quad d F=\eta \frac{d v}{d z} d S$
where_ $\eta$ is the dynamic viscosity coefficient depending on the nature of the fluid (and temperature).
The size of the viscosity coefficient is: $[\eta]_{S . I .}=\frac{\mathrm{N} \cdot \mathrm{m}}{\mathrm{m}^{2} \cdot \frac{\mathrm{~m}}{\mathrm{~s}}}=\frac{\mathrm{Kg} \cdot \mathrm{m} \cdot \mathrm{s}}{\mathrm{s}^{2} \cdot \mathrm{~m}^{2}}=\frac{\mathrm{Kg}}{\mathrm{m} \cdot \mathrm{s}}$ or $[\eta]_{\text {S.I. }}=\frac{\mathrm{N} \cdot \mathrm{S}}{\mathrm{m}^{2}}$
The CGS unit is one poise ( $\mathbf{P}$ ) (after POISEUILLE): $1 P=1 \frac{\mathrm{~g}}{\mathrm{~cm} \cdot \mathrm{~s}}=0,1 \frac{\mathrm{Kg}}{\mathrm{m} \cdot \mathrm{s}}$. Thus the unit of measurement in SI is DECAPOISE (daP).

For liquids $\eta$ is $10^{-3} \mathrm{daP}$ and for gases $\eta$ is two orders of magnitude lower: $\eta=10^{-5} \mathrm{daP}$.
KINEMATIC VISCOSITY $v$ is the ratio between viscosity $\eta$ and density $\rho$ of the fluid: $v=\frac{\eta}{\rho}$ and its unit of measure is $[v]_{S . I}=1 \frac{K g}{m \cdot s} \cdot \frac{m^{3}}{K g}=1 \frac{m^{2}}{s}$ in CGS the unit is STOKES (St) $1 S t=1 \frac{\mathrm{~cm}^{2}}{s}=10^{-4} \frac{\mathrm{~m}^{2}}{\mathrm{~s}}$.

For liquids $v \approx 10^{-6} \frac{m^{2}}{s}$, for gases $10^{-5} \frac{m^{2}}{s}$ (higher than for liquids)

The existence of internal friction (viscosity) can easily be shown experimentally by suspending a disk or a cylinder under which another disk or cylinder rotates by a wire.

The layer of fluid immediately adjacent to the rotating body adheres to it and is driven by it. The other layers are driven at lower and lower speeds until the last layer is glued to the other disc or cylinder, on which a frictional force will be exerted that will rotate it (see fig.1).


Figure 1

## POISEUILLE'S FORMULA

The study the stationary linear flow of a viscous fluid through a tube. Linear flow occurs at not too high speeds, or at not too large diameters.


Fig 2.

If we delimit a current tube of radius $r$ (see fig.2). The pressure forces from the extremities act on the fluid in this tube with the resultant $p_{1} \pi r^{2}-p_{2} \pi r^{2}$ internal friction force on the side surface, exerted by the rest of the fluid due to viscosity $\mathfrak{J} 2 \pi \mathrm{rl}$.

The flow being stationary (with constant speed ), the forces balance: $\left(p_{1}-p_{2}\right) \pi r^{2}=2 \pi r / \sqrt{s}$

$$
\begin{equation*}
\text { or }\left(p_{1}-p_{2}\right) r=-2 l \eta \cdot \frac{d v}{d z} \tag{3}
\end{equation*}
$$

( $\mathfrak{J}$ is the internal friction force: $\mathfrak{I}=\frac{d F}{d S}=\eta \frac{d v}{d z}$ ) The minus sign appears because of the minus sign of the velocity gradient $\frac{d v}{d z}<0$; the speed on the axis of the tube is maximum and decreases on the walls of the tube, being zero at the wall, in the adjacent layer.

By integration we get:

$$
\begin{align*}
& d v=-\frac{p_{1}-p_{2}}{2 l \eta} r d z  \tag{4}\\
& v=-\frac{p_{1}-p_{2}}{4 l \eta} r^{2}+C \frac{d v}{d z} \tag{5}
\end{align*}
$$

The constant C is determined from the condition that at the wall for $\mathrm{r}=\mathrm{R}$, the velocity should be zero.

$$
\begin{equation*}
v(r)=\frac{p_{1}-p_{2}}{4 l \eta}\left(R^{2}-r^{2}\right)=v_{m}\left(1-\frac{r^{2}}{R^{2}}\right) \tag{6}
\end{equation*}
$$

where $\quad v_{m}=\frac{\left(p_{1}-p_{2}\right) R^{2}}{4 l \eta}$
So $\quad \frac{v}{v_{m}}=1-\frac{r^{2}}{R^{2}} \quad$ - the speed distribution is parabolic.
If we calculate the volume flow rate:

$$
\begin{gather*}
Q_{v}=\int v d S=\int_{0}^{r} v \cdot 2 \pi r \cdot d r=\int_{0}^{r} \frac{2 \pi\left(p_{1}-p_{2}\right)}{4 l \eta}\left(R^{2}-r^{2}\right) r d r=\frac{2 \pi\left(p_{1}-p_{2}\right)}{4 l \eta} . \\
\cdot \int_{0}^{r}\left(R^{2}-r^{2}\right) r d r=\frac{\pi\left(p_{1}-p_{2}\right)}{8 l \eta} R^{4}=S v  \tag{9}\\
\text { where } \quad v=\frac{p_{1}-p_{2}}{8 l \eta} R^{2}
\end{gather*}
$$

$$
Q_{v}=\frac{\pi\left(p_{1}-p_{2}\right)}{8 l \eta} R^{4}=S v
$$

$$
\begin{equation*}
\text { where } \quad v=\frac{p_{1}-p_{2}}{8 l \eta} R^{2} \tag{10}
\end{equation*}
$$

This is POISEUILLE's formula (1841). The flow rate is proportional to the pressure drop per unit length of the tube and to the power of the 4th of the tube radius.

This formula can be used to determine the viscosity of fluids (for example in the OSTWALD viscometer).

Poiseuille's law explains some aspects of the physiology of blood circulation. The human capillary network totals $10^{5} \mathrm{Km}$. According to the needs of the body, the flow of blood is easily regulated by constriction or dilation of the blood vessels.

## BOUNDARY LAYER. STOKES' LAW .

When a body moves in a fluid, a very thin layer of fluid adheres to the surface, driven by the body. In a linear regime, that is, at not too high speeds, in the vicinity of the body there is a relatively thin layer, called the boundary layer, in which the speed drops to zero and in which the friction forces due to viscosity are manifested.

We can evaluate the thickness of the boundary layer as follows: we write the length and width of the boundary layer with $h$ and $b$. Then the force of internal friction according to Newton's law is: $F \approx \eta \cdot \frac{v}{d} \cdot h \cdot b$

On the other hand, the internal friction force can be found from the momentum variation of the fluid (from the layer with zero velocity to the one with velocity $v$ ) in the unit of time:

$$
\begin{equation*}
F \approx Q_{m} \cdot v=\rho d b v \cdot v=\rho d b \cdot v^{2} \tag{12}
\end{equation*}
$$

From these two expressions it follows: $d \cong \sqrt{\frac{h \eta}{\rho v}}$
If the ratio $\frac{h b}{d}=l$ is a characteristic length of the body, then the frictional force becomes:

$$
\begin{equation*}
F=\text { const } . \cdot \eta l v \text { Stokes' law } \tag{14}
\end{equation*}
$$

STOKES' LAW: The frictional force $F$ encountered by a body (in flow regime) is proportional to the viscosity of the fluid $\eta$, to the characteristic linear dimension of the body and to its speed $v$. In the case of the sphere, STOKES' formula $F=6 \pi \eta r v$ is obtained: (15) where r - is the radius of the sphere.

## EXPERIMENTAL ASSEMBLY

To determine the dynamic and kinematic viscosity, the Hoppler viscometer is used. The temperature can also be controlled from inside the viscometer in the range $\left(-20,+120^{\circ} \mathrm{C}\right)$


1. the tube where the ball falls
2. ball
3. the closing cap of the tube where the viscous fluid that will be studied is introduced
4. outer tube
5. thermometer
6. The temperature can be controlled inside the viscometer in the range $(-20 ;+120){ }^{\circ} \mathrm{C}$

## PROCEDURE FOR THE PRACTICAL WORK - STEPS

## MEASUREMENT OF BALL FALLING TIME

Record the time the ball falls between the two marks on the tube. The tube is filled in advance with the viscous liquid to be studied (the distance between the two marks has a value of $4.5 \mathrm{~cm}=45 \mathrm{~mm}$ )

## EXPERIMENTAL DATA PROCESSING

CALCULATION OF DYNAMIC VISCOSITY is done with the formula: $\eta=k \cdot\left(\rho_{1}-\rho_{2}\right) . t$
where
k is the constant of the ball (having the unit $\frac{m P a . s . c m^{3}}{g . s}$ )
$\rho_{1}$ - the density of the ball (having the unit: $\mathrm{g} / \mathrm{cm}^{3}$ )
$\rho_{2}$ - the density of the liquid to be studied (having the unit: $\mathrm{g} / \mathrm{cm}^{3}$ )
t - the ball's falling time.
CALCULATION OF KINEMATIC VISCOSITY is done with the formula: $v=\frac{\eta}{\rho_{2}}$
Unit of measure is $1 \frac{\mathrm{~mm}^{2}}{s}=1 \mathrm{cSt}$
The calculation of the density of the falling ball is determined as follows:
$\rho=\frac{m}{V}=\frac{m}{\frac{4 \cdot \pi \cdot R^{3}}{3}}=\frac{6 m}{\pi \cdot D^{3}}$
Where R- radius of the ball, D- diameter of the ball, m - mass of the ball.

Table 1：Material constants

| Article | Ball | Material | Density <br> $\mathrm{g} / \mathrm{cm}^{3}$ | Ball diameter <br> mm | Const K <br> $\mathrm{mPa} \cdot \mathrm{s} \cdot \mathrm{cm}^{3} / \mathrm{g} \cdot \mathrm{s}$ | The range of values <br> obtained after measuring <br> $\mathrm{mPa} \cdot \mathrm{s}$ |
| :--- | :--- | :--- | :--- | :--- | :--- | :--- |
| V91101 | 1 | Glass <br> Borosilicate | 2.2 | 15.81 | 0.007 | $0.6 \ldots 10$ |
| V91102 | 2 | Borosilicate glass | 2.2 | 15.6 | 0.09 | $7 \ldots 130$ |
| V91103 | 3 | Ni－Fe | 8.1 | 15.6 | 0.09 | $30 \ldots . .700$ |
| V91104 | 4 | Ni－Fe | 8.1 | 15.2 | 0.7 | $200 \ldots 4800$ |
| V91105 | 5 | Steel | $7.7-8.1$ | 14.0 | 4.5 | $800 \ldots 10000$ |
| V91106 | 6 | Steel | $7.7-8.1$ | 11.0 | 33 | $6000 \ldots 75000$ |

Table 2：Determination of dynamic and kinematic viscosity $\eta$（mPa．s）

| － | $\sum_{\sum}^{000}$ | 太气㐅⿸厂犬 |  | $\begin{aligned} & \overparen{O} \\ & \frac{\pi}{0} \\ & \frac{0}{0} \end{aligned}$ | $\stackrel{\cong}{\cong}$ | $\eta=k .\left(\rho_{1}-\rho_{2}\right) . t$ | $v=\frac{\eta}{\rho_{2}} \frac{\stackrel{\pi}{\mathrm{~s}}}{\text { a }}$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 1 | 4.58 | 1，580 | 0.01972 | 2.2212 |  |  |  |
| 2 | 4.41 | 1，560 | 0.12801 | 2.2199 |  |  |  |
| 3 | 16.13 | 1，559 | 0.13945 | 8.1380 |  |  |  |
| 4 | 14.91 | 1，518 | 0.78017 | 8.1480 |  |  |  |
| 5 | 11.16 | 1，399 | 6.32917 | 7.7925 |  |  |  |
| 6 | 5.42 | 1，099 | 31.15722 | 7.8147 |  |  |  |

Table 3: Determination of dynamic and kinematic viscosity

## for the plantusin syrup

| No. No. | (\%) | $\begin{aligned} & \rho_{2} \\ & \mathrm{~g} / \mathrm{cm}^{3} \end{aligned}$ | K $\frac{m P a . s . c m^{3}}{g . s}$ | $\rho_{1}$ $\left(\mathrm{g} / \mathrm{cm}^{3}\right)$ | t(s) | $\begin{aligned} & \eta=k .\left(\rho_{1}-\rho_{2}\right) . t \\ & \eta(m P a . s) \end{aligned}$ | $v=\frac{\eta}{\rho_{2}} \mathrm{~mm}^{2} / \mathrm{s}$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 1 | 20 | 1.06 | 0.12801 | 2.21996 |  |  |  |
| 2 | 40 | 1.1 |  |  |  |  |  |
| 3 | 60 | 1.14 |  |  |  |  |  |
| 4 | 80 | 1.18 |  |  |  |  |  |
| 5 | 100 | 1.24 |  |  |  |  |  |

## DYNAMIC, KINEMATIC AND RELATIVE VISCOSITY MEASUREMENT <br> UBELLOHDE VISCOSIMETER

## THEORETICAL INTRODUCTION

The forces of intermolecular attraction that oppose the relative displacement of neighboring molecules cause the appearance of the phenomenon of internal friction also called as viscosity The frictional forces between the layers of flowing liquid are given by
NEWTON'S FORMULA: $F=\eta \cdot S \cdot \frac{\Delta v}{\Delta x}$
Where $\eta$ is the viscosity coefficient,
S - contact surface between the liquid layers
$\frac{\Delta v}{\Delta x}$ - the speed gradient
POISEUILLE'S FORMULA is used to calculate the viscosity $\eta: Q=\frac{p \cdot R^{4} \cdot \pi}{8 \cdot \eta \cdot l}$
Q - flow; p-pressure; R- radius of the tube; 1-tube length.
The laminar flow of liquid is described by Poiseuille's equation: $Q=\frac{p \cdot R^{4} \cdot \pi}{8 \cdot \eta \cdot l}$
The definition of flow is: $Q=\frac{V}{t} \Rightarrow V=Q . t$
For the unknown liquid we have the following formula: $V=\frac{p \cdot R^{4} \pi \cdot t}{8 \eta l}$
For water: $V_{0}=\frac{p_{0} \cdot R^{4} \pi \cdot t_{0}}{8 \eta_{0} l}$
The pressure of the unknown liquid is: $p=\rho . g . h$
and the water pressure is: $p_{0}=\rho_{0} . g . h$
if the liquid has: $\mathrm{V}=\mathrm{V}_{0}$
$\frac{p \cdot R^{4} \pi \cdot t}{8 \eta l}=\frac{8 . \eta_{0} l}{p_{0} R^{4} \pi \cdot t_{0}}=1$
$\frac{p . t}{\eta}=\frac{p_{0} \cdot t_{0}}{\eta_{0}} \Rightarrow \frac{\rho . g . h . t}{\eta}=\frac{\rho_{0} g . h . t_{0}}{\eta_{0}}$
$\frac{\rho . t}{\eta}=\frac{\rho_{0} . t_{0}}{\eta_{0}}$
The viscosity formula is:
$\eta=\eta_{0} \frac{\rho . t}{\rho_{0} \cdot t_{0}}$
And the relative viscosity formula is:
$\eta_{\text {relatif }}=\frac{\eta}{\eta_{0}}=\frac{\rho . t}{\rho_{0} \cdot t_{0}}$

## EXPERIMENTAL ASSEMBLY

The figure below represents the Ubellohde viscometer . It is made of glass and has the following constructive parts:

- reservoir R1, R2, R3 and R4;
- communication tubes, open at the upper end, (A), (B) and (C).
- tube with controlled diameter (C) through which R3 and R4 reservoirs communicate.

The determination of the dynamic viscosity coefficient $(\eta)$ is based on the measurement of the flow time of a controlled volume of liquid through the tube (C). This process is governed by the
 pressure difference $(\Delta p)$ between the ends of the tube:

$$
\mathrm{Q}=\frac{\mathrm{V}}{\mathrm{t}}=\frac{\pi \cdot \mathrm{r}^{4} \cdot \Delta \mathrm{p}}{8 \cdot \mathrm{~L} \cdot \eta}(1)
$$

Figure 1. Ubbelohde viscometer

Practically, the flow time ( t ) of a liquid volume ( V ) is measured, and the dynamic viscosity coefficient is expressed by relation (2).

$$
\eta=\frac{\pi \cdot r^{r} \cdot \Delta p \cdot t}{8 \cdot L \cdot V}
$$

If the relation (2) is applied for a reference liquid for which the dynamic viscosity coefficient $\left(\eta_{r}\right)$ is known, the time $\left(\mathrm{t}_{\mathrm{r}}\right)$ is required to empty the volume $(\mathrm{V})$. In this case, the equation is written as follows: $\quad \eta_{r}=\frac{\pi \cdot r^{r} \cdot \Delta p \cdot t_{r}}{8 \cdot L \cdot V}$

By dividing member by member the relations (2) and (3), the following formula results: $\frac{\eta}{\eta_{r}}=\frac{t}{t_{r}}$ (4)
from where the coefficient of dynamic viscosity of the studied liquid $(\eta)$ is calculated according to the coefficient of dynamic viscosity of the reference liquid ( $\eta_{\mathrm{r}}$ ).
$\eta=\eta_{r} \cdot \frac{t}{t_{r}}=d \cdot \frac{t}{t_{r}}$ (5)

## PROCEDURE FOR THE PRACTICAL WORK - STEPS

1. For the practical execution of the measurement, fill the viscometer with the liquid to be studied through the communication tube (A) (a volume of approximately $14-15 \mathrm{ml}$ is recommended).
2. The communication tube $(\mathrm{C})$ is closed and with the help of the air pressure introduced in the tube (B) the liquid is pushed up in the tube (B) until it exceeds the reservoir R2. Both tubes (C) and (B) are released. Liquid begins to flow through tube (B) to reservoir (R4). When the liquid level reaches the narrow part between the tanks (R2) and (R3), the time counting starts (the timer starts). When the tank (R3) is empty, stop the timer and note the recorded time ( t ).
3. Empty the viscometer, wash and dry it. After that, filled it with the reference liquid for which the dynamic viscosity coefficient $\left(\eta_{\mathrm{r}}\right)$ is known and repeat the procedure described above. The emptying time ( $\mathrm{t}_{\mathrm{r}}$ ) of the reservoir ( R 3 ) is also noted this time.
4. The values (t) and ( $\mathrm{t}_{\mathrm{r}}$ ) are entered in relation (5) and the coefficient of dynamic viscosity $(\eta)$ of the studied liquid is calculated depending on that of the reference liquid $\left(\eta_{r}\right)$.

Table 1: Determination of relative viscosity for Glycerin

| No. | C (\%) | $d=\frac{\rho}{\rho_{0}}$ | t (s) | $\eta=d \cdot \frac{t}{t_{r}}$ | $\eta$ | $\mu$ | $\eta$ real |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 1 | water | 1 |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |
| 1 |  | 1.0253 |  |  |  |  |  |
| 2 | 10 |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |
| 1 | 20 | 1.0581 |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |
| 1 | 30 | 1.0851 |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |
| 1 | 40 |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |

Average values and errors will be calculated and introduced in the table

Table 2: Determination of the viscosity relative to the pine bud syrup depending on the concentration

| No. | C(\%) | $d=\frac{\rho}{\rho_{0}}$ | t(s) | $\eta_{\text {relativ }}=d \cdot \frac{t}{t_{0}}$ | $\bar{\eta}_{\text {relativ }}$ | $\mu$ | $\eta_{\text {rel real }}$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 1 | water | 1 |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |
| 1 | 10 |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |
| 1 | 20 |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |
| 1 | 30 |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |
| 1 | 40 |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |

Average values and errors will be calculated and introduced in the table

Table 3: Determination of the viscosity relative to Panadol syrup depending on the concentration

| No. | C(\%) | $d=\frac{\rho}{\rho_{0}}$ | t(s) | $\eta_{\text {relativ }}=d \cdot \frac{t}{t_{0}}$ | $\bar{\eta}_{\text {relativ }}$ | $\mu$ | $\eta_{\text {rel real }}$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 1 | water | 1 |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |
| 1 | 20 |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |
| 1 | 40 |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |
| 1 | 60 |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |
| 1 | 80 |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |
| 1 | 100 |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |

Average values and errors will be calculated and introduced in the table

Table 4: Determination of the relative viscosity for the plantusin solution


Average values and errors will be calculated and introduced in the table

Using the following formula from FRX:
$\eta=k \cdot t \cdot \rho$ - the dynamic viscosity using the Ubbelohde viscometer can be calculated Where $\mathrm{k}=0.005096 \frac{\mathrm{~m}^{2}}{\mathrm{~s}^{2}}$ the constant of the Ubbelohde viscometer
t is the flow time between the two landmarks
$\rho$ - the density of the liquid used in the experiment

Table 5: Determination of the dynamic viscosity of the plantusin solution

| No. | C(\%) | $\rho$ | t(s) | $\eta=k \cdot t \cdot \rho$ | $\eta$ | $\mu$ | $\eta_{\text {real }}$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 1 | water | 1 |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |
| 1 | 20 | 1.06 |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |
| 1 | 40 | 1.1 |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |
| 1 | 60 | 1.14 |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |
| 1 | 80 | 1.18 |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |
| 1 | 100 | 1.24 |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |

Average values and errors will be calculated and introduced in the table

Using the following formula from FRX:
$\boldsymbol{\eta}=\mathrm{k} \cdot \mathrm{t} \cdot \boldsymbol{\rho}$ - the dynamic viscosity using the Ubbelohde viscometer can be calculated Where $\mathrm{k}=0.005096$ - the constant of the Ubbelohde viscometer t is the flow time between the two landmarks
$\rho$ - the density of the liquid used in the experiment

Table 6: Determination of the kinematic viscosity for the plantusin solution

| No. | c(\%) | $\rho$ | t(s) | $v=\eta / \mathrm{c}$ |  | $\eta$ | $\mu$ | $\eta_{\text {real }}$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 1 | water | 1 |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |  |
| 1 | 20 | 1.06 |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |  |
| 1 | 40 | 1.1 |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |  |
| 1 | 60 | 1.14 |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |  |
| 1 | 80 | 1.18 |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |  |
| 1 | 100 | 1.24 |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |  |

Average values and errors will be calculated and introduced in the table

## THE LAWS OF THE GRAVITATIONAL PENDULUM

## THEORETICAL INTRODUCTION

The simple gravitational pendulum is a point suspended by an inextensible string of negligible mass, which can oscillate in a vertical plane, around a point of suspension under the action of its weight. Frictional forces are neglected.


Choosing for the deviation angle the positive trigonometric sense and the Oz axis perpendicular to the figure towards the reader, the force moment and the kinetic moment relative to the suspension point are written: $M_{z}=M=-m g l \sin \theta$

$$
\begin{aligned}
& L_{z}=L=m v l=m \dot{\theta} l^{2} \\
& M_{z}=\dot{L}_{z} \Rightarrow \ddot{\theta}=\frac{g}{l} \sin \theta=0 \\
& -m g l \sin \theta=m \ddot{\theta} l^{2} \\
& -\frac{g}{l} \sin \theta=\ddot{\theta} \Rightarrow \ddot{\theta}+\frac{g}{l} \sin \theta=0
\end{aligned}
$$

By simplifying the mass, it follows that the oscillations do not depend on the mass of the material point. The equation can be solved for oscillations under small angles $\theta \ll \operatorname{rad}\left(\theta<6^{0}\right) \Rightarrow \sin \theta \approx \theta$ in radians $\Rightarrow \ddot{\theta}+\frac{g}{l} \theta=0$.

$$
\ddot{\theta}+\omega^{2} \theta=0 \quad \omega_{\text {ha }} \overline{\bar{a} q} \sqrt{\frac{g}{q}} \Rightarrow \overline{\overline{i c}} \text { oscillator equation }
$$

The solution is $\theta=\alpha \cos (\omega t+\varphi)$
$\omega=\sqrt{\frac{g}{l}} \Rightarrow T=2 \pi \sqrt{\frac{l}{g}}$
The relationship for the period is obtained by simplifying the following expression:
$T=2 \pi \sqrt{\frac{l}{g}}\left[1+\left(\frac{1}{2}\right)^{2} \sin ^{2} \frac{\alpha}{2}+\left(\frac{1 * 3}{2 * 4}\right) \sin ^{4} \frac{\alpha}{2}+\cdots\right]$

## The laws of the gravitational pendulum:

- The law of substance or the law of masses: the period does not depend on the mass and the nature of the substance of the material point.
- The law of isochronism of small oscillations: small oscillations $\left(\theta \ll 6^{0}\right)$ do not depend on their angular amplitude.
- The law of lengths: the period of oscillations is directly proportional to the square root of the pendulum length and inversely proportional to the square root of the gravitational acceleration.


## EXPERIMENTAL ASSEMBLY

The apparatus consists of three gravity pendulums made of different materials (brass, lead, aluminum) and different masses, suspended from a support so that their lengths can be varied.

## PROCEDURE FOR THE PRACTICAL WORK - STEPS

- Law of masses: The three pendulums of the same length but of different masses are put into oscillation with small amplitude. The time of oscillation $(n \gg 10)$ is measured $n$, which divided by $n$ gives the period of oscillation.
- The law of lengths: it is studied with a pendulum having the same mass but different lengths. The three pendulums of the same mass but different lengths are set to oscillate with small amplitude in turn. The time of oscillation ( $n \gg 10$ ) is measured $n$, which divided by $n$ gives the period of oscillation.


## PROCESSING OF THE RESULTS

- The law of masses: Complete table 1. Calculate the errors using the formula and plot the period T according to the mass $\bar{T}=\bar{T}(m)$. With the method of the smallest, it is demonstrated that the period does not depend on the mass of the pendulum. where:

$$
\bar{T}=A m+B ; A=\frac{\overline{x y}-\bar{x} \bar{y}}{\overline{x^{2}}-(\bar{x})^{2}}, B=\frac{\overline{x^{2}} \bar{y}-\bar{x} \overline{x y}}{\overline{x^{2}}-(\bar{x})^{2}}, x=m, y=\bar{T}
$$

Table 1: Verifying the law of masses $1=$ $\qquad$ .cm


Mean squared error of the average: $\mu= \pm \sqrt{\frac{\sum_{i=1}^{n}\left(T_{i}-\bar{T}\right)^{2}}{n(n-1)}}$
$n$ - the number of measurements, $T_{\text {real }} \in(\bar{T}-\mu, \bar{T}+\mu), \bar{T}=\frac{\sum_{i=1}^{n} T_{i}}{n}$.

## LEAST SQUARES METHOD

There are laboratory works in which the aim is to determine a quantity that constitutes a constant coefficient from a linear relationship of the form: $\mathbf{Y}=\mathbf{A} \cdot \mathbf{x}+\mathbf{B}$.

Such a relation is Hooke's law, for example: $\mathrm{F}=\mathrm{kx}$, in which the determination of the coefficient k or the acceleration of the center of the cylinder rolling on an inclined plane with rolling friction is required. Depending on $\sin \alpha \approx \alpha$ ( for $\alpha-\operatorname{small}: a=A . \alpha+B$ ) where $B$ is required.

By making several measurements for x and y , the constants A and B can be determined from a graphic representation. Usually the representative points are not exactly on a straight line, being affected by accidental errors. The question arises by which methods the constants $A$ and $B$ should be determined so that their values are as close as possible to the real value.

This problem is solved by the "least squares method" which is based on the following condition: the most probable value of the constants will be the one for which the sum of the squared errors is minimal.

It can be written accordingly:

$$
\begin{aligned}
& \bar{x}=\frac{1}{n} \sum_{i} x_{i} \\
& \bar{y}=\frac{1}{n} \sum_{i} y_{i} \quad \begin{array}{l}
\text { So: } \overline{x^{2}} A+\bar{x} B-\overline{x y}=0 \\
\overline{x^{2}}=\frac{1}{n} \sum_{i} x_{i}^{2} \quad \bar{x} A+B-\bar{y}=0 \\
\overline{x y}=\frac{1}{n} \sum_{i} x_{i} y_{i}
\end{array}
\end{aligned}
$$

The solutions of these equations constitute the values sought for A and B

$$
A=\frac{\overline{x y}-\bar{x} \cdot \bar{y}}{\overline{x^{2}}-(\bar{x})^{2}}
$$

$$
B=\frac{\overline{x^{2}} \cdot \bar{y}-\bar{x} \cdot \overline{x y}}{\overline{x^{2}}-(\bar{x})^{2}}
$$

The method can also be applied in non-linear relationships obtained on the basis of similar demonstrations.

## The law of lengths:

Complete table 2.
The errors are calculated using the formula and the period T is represented graphically according to the length.
Table 2: Verifying the law of lengths $m=. . . . . . . . . . . . g$

|  | $l(\mathrm{~cm})$ | $t(s)$ | $T(s)$ | $\bar{T}(s)$ | $\mu$ | $T_{\text {real }}(s)$ | $T_{C}=2 \pi \sqrt{\frac{l}{g}}(s)$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 1 |  |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |
| 4 |  |  |  |  |  |  |  |
| 5 |  |  |  |  |  |  |  |
| 1 |  |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |
| 4 |  |  |  |  |  |  |  |
| 5 |  |  |  |  |  |  |  |
| 1 |  |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |
| 4 |  |  |  |  |  |  |  |
| 5 |  |  |  |  |  |  |  |
| 1 |  |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |
| 4 |  |  |  |  |  |  |  |
| 5 |  |  |  |  |  |  |  |

The mean squared error of the average: $\mu= \pm \sqrt{\frac{\sum_{i=1}^{n}\left(T_{i}-\bar{T}\right)^{2}}{n(n-1)}}$
$n$ - the number of measurements, $T_{\text {real }} \in(\bar{T}-\mu, \bar{T}+\mu), \bar{T}=\frac{\sum_{i=1}^{n} T_{i}}{n}$.
5. Conclusions: The period does not depend on the mass, it depends on the length of the pendulum.

## DETERMINATION OF THE MOMENT OF INERTIA

## OF A RIGID BODY USING THE TRIFILAR TORSION PENDULUM

## THEORETICAL INTRODUCTION

The trifilar torsion pendulum consists of two equal disks connected by three thin wires (of negligible mass), each of length " L ". The assembly is positioned vertically (Figure 1-1)

Let the radius of the discs be " $R$ ". The upper disc is fixed while the lower disc can rotate around the vertical axis "OQ". After rotating the lower disk around the "OQ" axis with the angle " $\alpha_{0}$ " compared to the initial equilibrium position (Fig 1-2), the movable disk rises from the initial position, so the potential energy of the system increases.

(1)

(2)

(3)

Fig. 1 Construction details of the trifilar torsion pendulum
Left free, the system returns to its initial position and engages in a circular oscillating motion. The oscillation of the mobile disk has a period " $\mathrm{T}_{\mathrm{o}}$ ", dependent on the length " L " of the supporting wires, the radius " R " of the disks, the initial angle of rotation " $\alpha_{0}$ " and the moment of inertia of the mobile disk, " $\mathrm{I}_{\mathrm{o}}$ ", in relation to the axis of rotation "OQ". During the oscillation the elongation [instantaneous value at time " t " of the angle, " $\alpha(\mathrm{t})$ "] changes periodically in time. The tension " F ", which acts in each supporting wire, is determined by the weight of the movable disc " $\mathrm{M}_{\mathrm{o}}$ ". The oscillation period is $T_{0}$. The relationship between these quantities is expressed by the differential equation (1).
$I_{0} \cdot \frac{d^{2} \alpha(t)}{\mathrm{dt}^{2}}+\frac{\mathrm{R}^{2} \cdot F(\mathrm{t})}{\mathrm{L}} \cdot \sin \alpha(\mathrm{t})=0$
In relation (1) the tension in the wire, " $\mathrm{F}(\mathrm{t})$ ", is dependent on the momentary value of the elongation " $\alpha(\mathrm{t})^{\prime}$, thus implicitly on time. In the case of small initial deviation $\left(\alpha_{0}<5^{\circ}\right)$, the acceptance of a few simplifications is justified:
(a) $\sin \alpha(t) \approx \alpha(t) \quad[\alpha(t)$ being expressed in radians]
(b) The tension in the wire is approximately constant and determined, at each moment of the movement, by the weight of the moving disk: $\mathrm{F} \approx \mathrm{M}_{\circ} \cdot \mathrm{g}$ ("g" is the gravitational acceleration).
(c) The oscillation is approximately aronic, that is, the period " $\mathrm{T}_{\mathrm{o}}$ " can be considered constant in time, independent of the angle " $\alpha_{0}$ ".
In the light of these simplifications, equation (1) can be transcribed as follows:

$$
\begin{equation*}
\frac{\mathrm{d}^{2} \alpha(\mathrm{t})}{\mathrm{dt}^{2}}+\frac{\mathrm{R}^{2} \cdot \mathrm{M}_{\mathrm{o}} \cdot \mathrm{~g}}{\mathrm{~L} \cdot \mathrm{I}_{\mathrm{o}}} \cdot \alpha(\mathrm{t})=0 \tag{2}
\end{equation*}
$$

A differential equation of the type

$$
\begin{equation*}
\frac{\mathrm{d}^{2} \alpha(\mathrm{t})}{\mathrm{dt}^{2}}+\omega_{0}^{2} \cdot \alpha(\mathrm{t})=0 \tag{3}
\end{equation*}
$$

has, as a solution, a function of the form: $\alpha(\mathrm{t})=\mathrm{A} \cdot \sin \alpha_{0} \cdot \mathrm{t}$ where " $\alpha_{0}$ " represents the pulsation, and " A " is the amplitude of the oscillation. The relationship between the pulsation " $\alpha_{0}$ " and the period " $\mathrm{T}_{\mathrm{o}}$ " of the oscillation is:

$$
\omega_{\mathrm{o}}=\frac{2 \cdot \pi}{\mathrm{~T}_{\mathrm{o}}}(4)
$$

Comparing equations (2), (3) and (4), it follows:

$$
\omega_{\mathrm{o}}^{2}=\frac{4 \cdot \pi^{2}}{\mathrm{~T}^{2}}=\frac{\mathrm{R}^{2} \cdot \mathrm{M}_{\mathrm{o}} \cdot \mathrm{~g}}{\mathrm{~L} \cdot \mathrm{I}_{\mathrm{o}}}(5)
$$

$\mathrm{I}_{\mathrm{O}}=\frac{\mathrm{T}_{\mathrm{O}}^{2} \cdot \mathrm{R}^{2} \cdot \mathrm{~g} \cdot \mathrm{M}_{\mathrm{o}}}{4 \cdot \pi \cdot \mathrm{~L}}$
Relation (6) expresses the moment of inertia " $I_{0}$ " of the mobile disc. If a rigid object (denoted by " P " in Figure 1-3) is deposited on the moving disk, having the moment of inertia " $I_{P}$ " relative to the axis of rotation "OQ", then the period of oscillation (" $\mathrm{T}_{0}$ ") is changed becomes "T"). In this case, the
total mass that generates the tension in the wire is $\mathrm{M}=\mathrm{M}_{\mathrm{o}}+\mathrm{M}_{\mathrm{P}}$, " $\mathrm{M}_{\mathrm{P}}$ " being the mass of the object "P". The total moment of inertia, "I" is the sum of the moments of inertia " $\mathrm{M}_{\mathrm{o}}$ " and " $\mathrm{MP}_{\mathrm{P}}$.

$$
\begin{equation*}
\mathrm{I}=\mathrm{I}_{\mathrm{o}}+\mathrm{I}_{\mathrm{P}}=\frac{\mathrm{T}^{2} \cdot \mathrm{R}^{2} \cdot \mathrm{~g}}{4 \cdot \pi \cdot \mathrm{~L}} \cdot\left(\mathrm{M}_{\mathrm{o}}+\mathrm{M}_{\mathrm{P}}\right) \tag{7}
\end{equation*}
$$

From this formula, the moment of inertia "IP" of the object "P" in relation to the axis of rotation "OQ" is obtained: $\mathrm{I}_{\mathrm{P}}=\frac{\mathrm{R}^{2} \cdot \mathrm{~g}}{4 \cdot \pi \cdot \mathrm{~L}} \cdot\left[\mathrm{~T}^{2} \cdot\left(\mathrm{M}_{\mathrm{o}}+\mathrm{M}_{\mathrm{P}}\right)-\mathrm{T}_{\mathrm{o}}^{2} \cdot \mathrm{M}_{\mathrm{o}}\right]$

## PROCEDURE FOR THE PRACTICAL WORK - STEPS

(1)The mass Mo of the movable disk is determined.
(2) The torsional oscillation of the pendulum is initiated by deflecting the mobile disk with a small angle ( $\alpha_{0}<5^{\circ}$ ).
(3) The duration " $\mathrm{T}_{\mathrm{N}}$ " of a number of " N " complete oscillations is determined. For greater precision, choose as many complete oscillations as possible.If the oscillations amortize quickly, repeat the experience several times.
(4)The " $\mathrm{T}_{\mathrm{o}}$ " period is calculated by dividing the " $\mathrm{T}_{\mathrm{N}}$ " period by the number " N ".
(5) The mass " Mp " of the studied body is determined.
(6) The studied body settles in the center of the mobile disc.
(7) System oscillation is initiated ( $\left.\alpha_{0}<5^{\circ}\right)$.
(8) Follow the duration of " N " complete oscillations and obtain the period " T " by dividing the duration by the number of " N " complete oscillations.
(9) The data obtained are replaced in relation (8) to obtain the value of the moment of inertia "Ip" of the object in relation to the axis of rotation "OQ" of the system.

The moment of inertia in relation to parallel axes of rotation (verification of the Huygens-Steiner theorem )

Huygens - Steiner theorem refers to the moment of inertia of a body in relation to parallel axes of rotation (Figure 2). If the body has the moment of inertia "I $\mathrm{I}_{\mathrm{G}}$ " with respect to an axis of rotation passing through the center of gravity of the body ("G"), then the moment of inertia of the body ("Ip")
with respect to an axis of rotation parallel to the first axis, which passes through the point "P" located at the distance " d " from the center of gravity " G ", is

$$
\begin{equation*}
\mathrm{I}_{\mathrm{P}}=\mathrm{I}_{\mathrm{G}}+\mathrm{M} \cdot \mathrm{~d}^{2} \tag{9}
\end{equation*}
$$

where " M " represents the mass of the body.


Fig. 2 The essence of the Huygens - Steiner theorem
Huygens-Steiner theorem can be verified if two identical bodies, each of mass " m ", are placed on the moving disk, symmetrically in the hood with the center of the moving disk, the distance between the two bodies being " $2 \cdot \mathrm{~d}$ ". In this case, the moment of inertia of the assembly thus formed is with the relation:
$\mathrm{I}=\frac{\mathrm{R}^{2} \cdot \mathrm{~g}}{4 \cdot \pi \cdot \mathrm{~L}} \cdot\left[\mathrm{~T}^{2} \cdot\left(\mathrm{M}_{\mathrm{o}}+2 \cdot \mathrm{M}_{\mathrm{p}}\right)-\mathrm{T}_{\mathrm{o}}^{2} \cdot \mathrm{M}_{\mathrm{o}}\right]$
It is checked if the moment of inertia "I", thus determined, satisfies the relation (11).

$$
\begin{equation*}
\frac{\mathrm{I}}{2}=\mathrm{I}_{\mathrm{o}}+\mathrm{m} \cdot\left(\frac{\mathrm{~d}}{2}\right)^{2} \tag{11}
\end{equation*}
$$

## EXPERIMENTAL DATA PROCESSING

To determine the moments of inertia with the above relationships, the following tables of values will be drawn up:

Table 1: Determination of the moment of inertia with the help

| No. <br> crt. | $\mathrm{M}_{\mathrm{o}}(\mathrm{kg})$ | $\mathrm{R}(\mathrm{cm})$ | $\mathrm{L}(\mathrm{cm})$ | $\mathrm{m}_{2}(\mathrm{~kg})$ | to(s) | N <br> $(\mathrm{Osk})$ | $\mathrm{T}_{\mathrm{o}}(\mathrm{s})$ | Yo <br> $\left(\mathrm{kg} . \mathrm{m}^{2}\right)$ |
| :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- |
| 1 |  |  |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |  |
| 4 |  |  |  |  |  |  |  |  |
| 5 |  |  |  |  |  |  |  |  |

Table 2:

| No. <br> crt. | m <br> $(\mathrm{kg})$ | $\mathrm{R}(\mathrm{cm})$ | $\mathrm{r}(\mathrm{cm})$ | $\mathrm{l}(\mathrm{cm})$ | m <br> $(\mathrm{kg})$ | $\mathrm{t}_{( }(\mathrm{s})$ | N <br> $($ Osk $)$ | $\mathrm{T}(\mathrm{s})$ | $\mathrm{I}\left(\mathrm{kg} . \mathrm{m}^{2}\right)$ |
| :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- |
| 1 |  |  |  |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |  |  |
| 4 |  |  |  |  |  |  |  |  |  |
| 5 |  |  |  |  |  |  |  |  |  |

## 5. Conclusions:

conclusions are drawn regarding the values of the moments of inertia thus obtained and their values

## THE STUDY OF BERNOULLI'S LAW

## THEORETICAL INTRODUCTION

## BERNOULLI'S EQUATION

## An incompressible fluid without viscosity is called IDEAL FLUID.

With fairly good accuracy, ideal liquids can be considered incompressible. Although gases are easily compressible, their compressibility is manifested only in flows with speeds close to the speed of sound. Experience shows that a large number of liquids and even gases, if their velocities are much lower than the speed of sound and the velocity gradients are small, correspond satisfactorily to the ideal fluid. We will study in this paragraph the flow of an ideal fluid.


Figure 1: Helpful schematic representation for Bernoulli's law

Let's consider the ideal fluid and an infinitely thin current tube, delimited at the ends by two sections $\mathrm{dS}_{1}, \mathrm{dS}_{2}$ and apply the fluid quantities thus delimited by the mechanical energy variation theorem. During dt the system moves along the tube from position AB to A'B'. Everything happens as if the mass disappears from the portion $\mathrm{AA}^{\prime}$

$$
\begin{equation*}
d m=\rho d S_{1} v_{1} d l=\rho d V \tag{1}
\end{equation*}
$$

with the kinetic energy $\frac{d m \cdot v_{1}{ }^{2}}{2}$ and the potential energy $\mathrm{dm} \cdot g \cdot h_{1}$, and in the portion BB' there would appear a mass equal (due to the incompressibility of the ideal fluid) $d m=\rho d S_{2} \cdot v_{2} \cdot d t=\rho d V$ with the kinetic energy $\frac{d m \cdot v_{2}{ }^{2}}{2}$ and the potential energy $d m \cdot g \cdot h_{2}$. The variation of the mechanical energy (kinetic and potential) of the system is therefore:

$$
\begin{equation*}
d E=\frac{1}{2} d m\left(v_{2}^{2}-v_{1}^{2}\right)+d m \cdot g\left(h_{2}-h_{1}\right) \tag{2}
\end{equation*}
$$

and must be equal to the mechanical work of the pressure forces exerted on the considered system:

$$
\begin{align*}
& d w=F_{1} \cdot d S_{1}-F_{2} \cdot d S_{2} \\
& =p_{1} d S_{1} v_{1} d t-p_{2} d S_{2} v_{2} d t  \tag{3}\\
& =p_{1} d V-p_{2} d V
\end{align*}
$$

and $\quad d V=\frac{d m}{\rho}$
The lateral pressure forces do not perform mechanical work and we have no tangential frictional forces, the fluid being assumed to be ideal (without viscosity) $\mathrm{dE}=\mathrm{dw}$

$$
\begin{equation*}
\text { So: } p_{1}+\frac{1}{2} \rho v_{1}^{2}+\rho g h_{1}=p_{2}+\frac{1}{2} \rho v_{2}^{2}+\rho g h_{2} \tag{4}
\end{equation*}
$$

or

$$
\begin{equation*}
p+\frac{1}{2} \rho v^{2}+\rho g h=\text { const } . \tag{5}
\end{equation*}
$$

(5) is the BERNOULLI equation

The constant generally differs from one current line to another (it is the same for flow without eddies).
The pressure p is the static pressure
$p d=\frac{1}{2} \rho v^{2}$-dynamic pressure (due to the kinetic energy of the fluid, being equal to the kinetic energy of the volume unit)
$\underline{p=\rho g h}$-positional pressure (or potential due to the potential energy being equal to the potential energy of the volume unit)

## BERNOULLI'S LAW

The total pressure in a fluid is constant along a streamline .
Static pressure is exerted on a surface element placed parallel to the current lines, for example on the walls of the tube (that's why it is sometimes called pressure on the walls) or exerted on an arbitrarily oriented surface element but moved in solidarity with the fluid. The static pressure is measured with the pressure probe.
The dynamic pressure is measured with the PRANDTL tube.
The total pressure is exerted on a surface element placed perpendicular to the current lines and is measured with the PITOT tube having the opening placed perpendicular to the current lines.

For the horizontal pipe, the BERNOULLI equation becomes:

$$
\begin{equation*}
p_{1}+\frac{1}{2} \rho v_{1}^{2}=p_{2}+\frac{1}{2} \rho v_{2}^{2}=p+\frac{1}{2} \rho v^{2}=\text { const } . \tag{6}
\end{equation*}
$$

## THE HYDRODYNAMIC PARADOX

Entering a narrow portion of a current tube, the speed of the particles increases (the flow rate being constant), that is, they move faster, from which it follows that the pressure in the fluid in the wide portion of the tube must be higher than in the
 narrow portion.

Figure 2
Indeed, because in the narrow portions of the tube the speed of the fluid increases (according to the continuity equation), the dynamic pressure also increases, therefore the static pressure, i.e. the pressure on the walls, must decrease so that their sum remains constant, according to Bernoulli's equation (HYDRODYNAMIC PARADOX )


Fig3:
the phenomenon of suction can appear, on which the name of the practical application is based, such as SPRAYER, WATER PIPE, BUNSEN LAMP

## EXPERIMENTAL ASSEMBLY

The device for verifying Bernoulli's law is a device Type: DVLB-1
The experiences that can be carried out with the help of this device, as well as the relative errors, listed in parentheses, are the following:
1.Measurement of total pressure using a pitot tube (2\%);
2. Partial pressure measurement ( $2 \%$ );
3. Verification of Bernoulli's law, total pressure =ct . $(2 \%)$;
4. Highlighting the turbulent flow;
5. Determination of dynamic pressure using the Prandtl tube. 3\%);
6.Measuring the flow rate of a fluid through a pipe (3\%);
7. Verification of the law of continuity, $\mathrm{Sv}=\mathrm{ct}(3 \%)$;
8. Highlighting the ascension force-qualitative.


Fig.1.Schema principială a dispozitivului de verificare a legii Iui Bernoulli.


Fig.2.Vedere de ansamblu a dispozitivului de verificare a legii lui Bernoulli. Vedere frontală.


Fig.4.Vedere de ansamblu a aranjării elementelor suplimentare la dispozitivul de verificare a legii lui Bernoulli pentru evidentierea fortei ascensională.

Apart from the verification of Bernoulli's law, the apparatus allows the experimental realization of other laws and phenomena. This device can also be used as a Pitot tube or a Prandtl tube. It is designed to be used together with a regular vacuum cleaner.

## DESCRIPTION

Fig. 1 shows very simply the constructive principle of operation of the Bernoulli's law verification device. In a tube, TB , of variable section, from 20 mm to 16.5 mm in diameter, a jet of air is allowed to circulate, which is sucked in the direction of the arrow in fig. 1 , with the help of a vacuum cleaner, which is not shown in this figure and which connects to the right end, as we look at the image.

The device is equipped with two pairs of pressure measuring probes, as follows: 1 and 3 are the total pressure measuring probes, and 2 and 4 are the static (wall) pressure measuring probes. On the two sections, the two pairs of probes are arranged in a cross, as can be seen in fig.2. On the front side, PF , of the device there is a measuring ruler in front of which a mercury manometer, MU , is positioned.

The mercury manometer, MU, will be used to measure static, dynamic, total pressures relative to atmospheric pressure or differentially. The coupling of this manometer, to the two pairs of probes, is done by means of polyethylene tubes. At the end 5 visible on the front panel PF , the coupling is made by means of a rubber sleeve, MC , and at the second end, 6 , the one that enters through the front panel, the coupling is made through the back, the coupling piece being provided with a spigot appropriate output.

## PROCEDURE FOR THE PRACTICAL WORK - STEPS

Mounting the device is as simple as possible, $90 \%$ is already done by the manufacturing company. The corresponding steps are as follows:

1- carefully read the instructions, as well as the CD with the experimental procedures;
2-remove the manometric tube, from the appropriate packaging, remove the safety rubber sleeves during transport, from the two ends. Unwrapping them is done starting with the longer end and then the shorter;

3-at the short end, 5 , connect a rubber sleeve, MC ;
4-insert the end 6 of the mercury manometer, MC , in the upper part of the front panel, as shown in fig.2. The insertion is done carefully, with a rotation of $30^{0}$ to the right and left of the vertical position. A rotation at a greater angle is forbidden, there is a risk that the mercury will come out of the capillary tube of the manometer.

5 -moderately tighten the M5 screws on the back of the device. Tightening is carried out until the capillary tube can no longer be rotated in the rubber gasket, from the device support.

The operations from 2-5 are done, either with the pressure tube inserted $75 \%$ in a plastic bag or on top of a bowl of water. These precautions are mandatory in order to avoid accidental spillage of mercury in the laboratory.

6- the capillary tube is inserted into a small transparent plastic bag, bag delivered together with the device and which is attached with adhesive tape to the front of the device. The role of this bag is to avoid spilling the mercury in the laboratory when the capillary tube accidentally breaks.

## VERIFICATION OF BERNOULLI'S LAW, $\mathrm{P}_{\mathrm{T}}=\mathrm{CT}$.

## The calculation formulas are:

$p_{1}+\frac{1}{2} \rho v_{1}{ }^{2}+\rho g h_{1}=p_{2}+\frac{1}{2} \rho v_{2}{ }^{2}+\rho g h_{2}$
If h1=h2
$p_{1}+\frac{1}{2} \rho v_{1}^{2}=p_{2}+\frac{1}{2} \rho v_{2}^{2}=p+\frac{1}{2} \rho v^{2}=$ const .
$\Delta p=p_{1}-p_{2}=\frac{\rho}{2}\left(v_{1}^{2}-v_{2}^{2}\right)=\frac{\rho}{2} \cdot v_{1}^{2} \cdot\left[\left(\frac{S_{1}}{S_{2}}\right)^{2}-1\right]$
At the same time $\Delta p=\rho_{\text {lichid }} . g . h$
So the flow speed of the liquid in the pipe is: $v_{1}=\sqrt{\frac{2 \rho_{\text {lichid }} \cdot g \cdot h}{\rho_{\text {aer }}\left[\left(\frac{S_{1}}{S_{2}}\right)^{2}-1\right]}}$
And the surface of the pipe is: $S_{1}=\frac{\pi \cdot D_{1}^{2}}{4}$ respectively: $S_{2}=\frac{\pi \cdot D_{2}^{2}}{4}$
So the speed of the fluid flow through the pipe is:
$v_{1}=\sqrt{\frac{2 \rho_{\text {lichid }} \text {.g.h }}{\rho_{\text {aer }}\left[\left(\frac{D_{1}}{D_{2}}\right)^{4}-1\right]}}$
And the speed in section $\mathrm{S}_{2}$ is: $v_{2}=v_{1} \frac{S_{1}}{S_{2}}=v_{1}\left(\frac{D_{1}}{D_{2}}\right)^{2}$
The data are entered in table 1. for the following cases:

1. Small closure with closed valves
2. Small closure with open valves
3. InLarge shutoff with open valves
4. Large closure with closed valves
5. With total opening

In tables $2,3,4$, the values of the pressure variation, of the flow speed through the pipe v 1 and v 2 are entered.

## 6. Experimental data processing

Table 1.

| No. crt | d1 (mm) | $\begin{aligned} & \hline \mathrm{d} 2 \\ & \mathrm{~mm} \end{aligned}$ | $\rho_{\text {air }} \mathrm{kg} / \mathrm{m}^{3}$ | $\rho$ liquid ( $\mathrm{kg} / \mathrm{m}^{3}$ ) | h (mm) | $\Delta \mathrm{p}\left(\mathrm{N} / \mathrm{m}^{2}\right)$ | v1 (m/s ${ }^{\mathbf{2}}$ ) | v2 (m/s ${ }^{2}$ ) |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 1 | 20 | 18 | 1.20 | 800 <br> Technical <br> alcohol |  |  |  |  |
| 2 |  |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |  |
| 4 |  |  |  |  |  |  |  |  |
| 5 |  |  |  |  |  |  |  |  |

Table 2 Pressure Calculation

|  | $\Delta p\left(\frac{N}{m^{2}}\right)$ | $\Delta \bar{p}\left(\frac{N}{m^{2}}\right)$ | $\mu_{\Delta p}\left(\frac{N}{m^{2}}\right)$ | $\Delta p_{\text {real }}\left(\frac{N}{m^{2}}\right)$ | $\Delta p=\rho_{\text {lichid }} \cdot g \cdot h$ <br> Average value: |
| :---: | :---: | :---: | :---: | :---: | :---: |
| 1 |  |  |  |  | $\sum \Delta p_{i}$ |
| 2 |  |  |  |  | $\Delta \bar{p}=\frac{i=1}{n}$ |
| 3 |  |  |  |  | Error: |
| 4 |  |  |  |  | $\sqrt{\sum^{n}\left(\Delta p_{i}-\Delta \bar{p}\right)^{2}}$ |
| 5 |  |  |  |  | $\mu_{\Delta p}= \pm \sqrt{\frac{i=1}{n(n-1)}}$ <br> Range of actual values: $\Delta p_{\text {real }} \in\left(\Delta \bar{p}-\mu_{\Delta p}, \Delta \bar{p}+\mu_{\Delta p}\right)$ |

Table 3: calculation of flow velocity v1

| No. crt | $v_{1}\left(\frac{m}{s}\right)$ | $\bar{v}_{1}\left(\frac{m}{s}\right)$ | $\mu_{v_{1}}\left(\frac{m}{s}\right)$ | $v_{1 \text { real }}\left(\frac{m}{s}\right)$ | $v_{1}=\sqrt{\frac{2 \rho_{\text {lichid }} \cdot g . h}{\left[\left(D_{1}\right)^{4},\right]}}$ |
| :---: | :---: | :---: | :---: | :---: | :---: |
| 1 |  |  |  |  | $\sqrt{\rho_{a e r}\left[\left(\frac{D_{1}}{D_{2}}\right)-1\right]}$ |
| 2 |  |  |  |  | Average value: |
| 3 |  |  |  |  | $\sum^{n} v_{1 i}$ |
| 4 |  |  |  |  | $\bar{v}_{1}=\frac{i=1}{n}$ |
| 5 |  |  |  |  | Error: $\mu_{v 1}= \pm \sqrt{\frac{\sum_{i=1}^{n}\left(v_{1 i}-\bar{v}_{1}\right)^{2}}{n(n-1)}}$ <br> Range of actual values: $v 1_{\text {real }} \in\left(\bar{v}_{1}-\mu_{v 1}, \bar{v}_{1}+\mu_{v 1}\right)$ |

Table 4: calculation of the flow speed $v_{2}$

| No. crt | $v_{2}\left(\frac{m}{s}\right)$ | $\bar{v} 2\left(\frac{m}{s}\right)$ | $\mu \nu 2\left(\frac{m}{s}\right)$ | $\nu_{2 \text { real }}\left(\frac{m}{s}\right)$ | $v_{2}=v_{1} \frac{S_{1}}{S_{2}}=v_{1}\left(\frac{D_{!}}{D_{2}}\right)^{2}$ <br> Average value: $\bar{v} 2=\frac{\sum_{i=1}^{n} v_{2 i}}{n}$ <br> Error: $\mu_{v 2}= \pm \sqrt{\frac{\sum_{i=1}^{n}\left(v 2_{i}-\bar{v} 2\right)^{2}}{n(n-1)}}$ <br> Range of actual values: $v 2_{\text {real }} \in\left(\bar{v} 2-\mu_{v 2}, \bar{v} 2+\mu_{v 2}\right)$ |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 1 |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |
| 4 |  |  |  |  |  |  |
| 5 |  |  |  |  |  |  |

## STUDY OF GAS LAW

## BOYLE-MARIOTTE LAW

## INTRODUCTION

The Boyle-Mariotte law (also called Boyle's law by Anglophones, or Mariotte's law or BoyleMariotte law by Francophones) is one of the thermodynamic laws of real gases

Mariotte and Boyle found with a difference of a few percent that the curve was close to an equilateral hyperbola given by Clapeyron 's coordinates ( $\mathrm{p}, \mathrm{V}$ ), that is:
$\mathrm{pV}=$ constant for a given constant temperature
Or it can be written $p_{1} \cdot V_{1}=p_{2} \cdot V_{2}$
That is, maintaining a constant temperature during an increase in pressure, there is a decrease in volume

The law was discovered by the Irishman Robert Boyle in 1662 and by the French Français Edme Mariotte in 1676 .

Guillaume Amontons stated in 1702 that this law is valid only at constant temperature and low pressure.

This relationship between pressure and volume was noted by 2 scientists Richard Towneley and Henry Power. Robert Boyle confirmed and experimentally discovered this law. Robert Hook built the apparatus for the experiment. Boyle's law is based on experiments carried out with air, which he considers to be composed of fluid particles at rest between two invisible springs.
Boyle was interested in air as an important element of life.
Boyle used a J-shaped tube in the experiment where he introduced mercury into a branch and forced the air to contract under the pressure of the mercury. He performed the experiment several times under different conditions: different amounts of mercury in the installation and concluded that the gas pressure varies inversely proportional to the volume of the occupied gas. The French physicist Edme Mariotte (1620-1684) discovered

The same law with Boyle independently of him in 1676, but Boyle published the results already in 1662. However, this law is called the Boyle-Mariotte law

Later in 1867 in Philosophiæ Naturalis Principia Mathematica, Newton

He demonstrated mathematically that an elastic fluid contains particles at rest, that repulsive forces are exerted between them that are inversely proportional to the distance

Density is inversely proportional to pressure, but this mathematical reality is not the physical explanation for the observed relationship.
The kinetic-molecular theory introduced by Maxwell and Boltzmann two centuries later is needed to explain this law.

This law is the first physical law expressed in the form of an equation that describes the dependence of two variable quantities.

## The mathematical equation for Boyle's law is the following:

$$
\mathrm{PV}=\mathrm{k} \text { or } \mathrm{P} 1 \mathrm{~V} 1=\mathrm{p} 2 \mathrm{~V} 2=\ldots=\mathrm{pnVn}
$$

P represents the pressure in the system
V represents the volume of the gas
k is a constant that depends on the pressure and volume in the system

## TRANSFORMATION (T=CONST.) - BOYLE-MARIOTTE LAW

$\mathrm{PV}=\mathrm{k}$ or $\mathrm{P} 1 \mathrm{~V} 1=\mathrm{p} 2 \mathrm{~V} 2=\ldots=\mathrm{pnVn}$
$(A) \rightarrow(B)$ Isothermal expansion (with $T_{1}=$ const )
(C) $\rightarrow$ (D) Isothermal compression ( $\mathrm{T}_{2}=$ const) and $\mathrm{T}_{1}>\mathrm{T}_{2}$




## PART I: VERIFICATION OF BOYLE MARIOTTE'S GAS LAW

## EXPERIMENTAL DEVICE

The instrument consists of burette 1 where the gas volumes $\mathrm{V} 1, \mathrm{~V} 2, \ldots, \mathrm{Vn}$ can be determined. By modifying the height of the reservoir 3, the volume of the gas in the burette changes, but in the same time the pressure also changes, which can be read on the manometer 2.


Figure 1. Schematic representation of the device used to verify Boyle Mariotte's gas law

## PROCEDURE FOR THE PRACTICAL WORK - STEPS

1. determine the pressure to be studied using the manometer 2 . You can write $p_{1}, p_{2}, \ldots p_{n}$
2. determine the corresponding volume from burette 1 . It can be written: $V_{1}, V_{2}, V_{3}, \ldots V_{n}$
3. 5 measurements are made for $p$, respectively for $V$
4.5 appropriate calculations are made to check $p_{1} V_{1}=p_{2} V_{2}=\ldots=p_{n} V_{n}$
4. Introduce the experimental data in the table below

Table 1. Gas law verification

| No. crt. | $\begin{array}{\|l\|} \hline V_{\text {read }} \\ \mathbf{c m}^{3} \end{array}$ | $\begin{aligned} & \mathbf{V}=\mathbf{V}_{\text {read }}+ \\ & \mathbf{V}_{\mathbf{0}} \\ & \mathbf{c m}^{3} \end{aligned}$ | $\begin{aligned} & \hline \mathbf{P}_{\text {read }} \\ & \mathrm{mm} \\ & \mathrm{H}_{2} \mathrm{O} \end{aligned}$ | $\begin{aligned} & \mathbf{P}=\mathbf{P}_{\text {read }}+\mathbf{p}_{0} \\ & \mathbf{m m H}_{2} \mathrm{O} \end{aligned}$ | $\begin{aligned} & \mathbf{p \mathbf { V } = \mathbf { K }} \\ & \frac{N}{m^{2}} \cdot m^{3}= \end{aligned}$ | $\begin{array}{r} \overline{\boldsymbol{K}} \\ N . m \end{array}$ | $\mu= \pm \sqrt{\frac{\sum\left(\kappa_{i}-\bar{K}\right)^{2}}{n(n-1)}}$ | $K_{\text {real }}$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 1 |  |  |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |  |
| 4 |  |  |  |  |  |  |  |  |
| 5 |  |  |  |  |  |  |  |  |

where
$\mathrm{P}_{\mathbf{0}}=\mathbf{7 6 0} \mathbf{~ m m ~ H g}=\mathbf{1 0 3 3 6} \mathbf{~ m m ~ H} \mathbf{2} \mathbf{O}=\mathbf{1 0 1 3 2 5} \mathbf{N} / \mathbf{m}^{2}-$ is the normal atmospheric pressure
And $\mathbf{V}_{\mathbf{0}}=\mathbf{1 2 3 . 8 8} \mathbf{~ c m}^{\mathbf{3}}$ is the constant of the device or it is also called "dead volume" - it is the volume of the manometer with liquid measured from division 0 to the liquid in the manometer.

To determine the volume $\mathrm{V}_{0}$, the isothermal transformation formula is used:
$\mathrm{pV}=$ constant as follows: $p_{1} \cdot\left(V_{1}+V_{0}\right)=p_{2}\left(V_{2}+V_{0}\right)$
$V_{0}=\frac{p_{2} V_{2}-p_{1} V_{1}}{p_{1}-p_{2}}$
$\begin{aligned} &{ }^{2} \begin{array}{l}p_{1}\end{array}=p_{0}+\rho . g . h_{1} \\ & p_{2}=p_{0}+\rho . g . h_{2}\end{aligned}$
After making the calculations, we obtain: $\mathrm{V}_{0}=123.88 \mathrm{~cm}^{3}=123.88 \mathrm{~mL}$
The following conversions can be made:
for pressure p: $1 \mathrm{mmH}_{2} \mathrm{O}=9.8 \mathrm{~N} / \mathrm{m}^{2}$
for volume $\mathrm{V}: 1 \mathrm{~cm}^{3}=1 \cdot 10^{-6} \mathrm{~m}^{3}$

Average values, errors and interval of real values can be calculated:
Average value for $\mathrm{pV}=\mathrm{K}: \overline{\mathrm{K}}=\frac{\sum_{i=1}^{n} K_{i}}{n}$ and
The mean squared error of the average: $\mu= \pm \sqrt{\frac{\sum\left(K_{i}-\bar{K}\right)^{2}}{n(n-1)}}$

## FINAL CONCLUSIONS

## Errors can be calculated for the values obtained by direct measurements of $\mathbf{p}$ and $\mathbf{V}$

The graphs can be represented by respecting the function: $p=p(V)$

## PART 2: CALCULATION OF THE UNIVERSAL CONSTANT OF GASES (R):

The universal perfect gas constant denoted by R is the product of Avogadro's number Na and Boltzman's constant kB

This product is $\mathbf{R}=\mathbf{8 . 3 1 4 4 6 2 1} \mathbf{J} \cdot \mathbf{~ m o l}^{-1} \cdot \mathbf{K}^{-1}$ with an accuracy of $7.5 \times 10^{-6} \mathrm{~J} \cdot \mathrm{~mol}^{-1} \cdot \mathrm{~K}^{-11}$

In the CGS system it can also have the following values:
$\mathrm{R}=0.0821 \mathrm{~L} \cdot \mathrm{~atm} /(\mathrm{mol} . \mathrm{K})$
$\mathrm{R}=8.2057 \times 10^{-5} \mathrm{~m}^{3} \mathrm{~atm} /(\mathrm{mol} \mathrm{K})$
$\mathrm{R}=62.3637 \mathrm{~L} \cdot \operatorname{Torr} /(\mathrm{mol} \cdot \mathrm{K})$ or $\mathrm{L} \cdot \mathrm{mmHg} /(\mathrm{mol} \cdot \mathrm{K})^{2}$
$\mathrm{R}=1.987 \mathrm{cal} \mathrm{K}^{-1} \mathrm{~mol}^{-1} \underline{3}$
Conversions of measurement units:

$$
\begin{aligned}
& 1 \mathrm{~Pa} \simeq 7,500615 \times 10^{-3} \mathrm{mmHg}(\text { millimeters of mercury column }) \\
& 1 \mathrm{~Pa} \simeq 9.869233 \times 10^{-6} \mathrm{~atm}(\text { normal atmosphere }) \\
& 1 \mathrm{~Pa}=10^{-5} \mathrm{bar} \\
& 1 \mathrm{hPa}=10^{-3} \mathrm{bar}=1 \text { mbar or millibar } \\
& 1 \mathrm{~Pa} \simeq 1.42 \times 10^{-4} \mathrm{PSI} \\
& 1 \mathrm{~Pa}=1 \mathrm{~N} / \mathrm{m}^{2}
\end{aligned}
$$

either: $10^{5} \mathrm{~Pa}=1 \mathrm{bar}=1000 \mathrm{mbar}=10.2 \mathrm{~m} \mathrm{H}_{2} \mathrm{O}$ column $($ water column meters $)=0.987 \mathrm{~atm}$
$1 \mathrm{psi}=6894.76 \mathrm{~Pa}\left(\mathrm{~N} / \mathrm{m}^{2}\right)=0.069 \mathrm{bar}=0.069 \mathrm{kgf} / \mathrm{cm}^{2}=51.7 \mathrm{mmHg}=689.7 \mathrm{~mm} \mathrm{H}_{2} \mathrm{O}=0.068$ atmospheres
By using the equation of state: $\mathrm{pV}=\mathrm{nRT}$
$R=\frac{p \cdot V}{n \cdot T}$
Where: $\mathrm{T}_{0}=273 \mathrm{~K}$ and the molar volume of air is: $\mathrm{V}=22.4 \mathrm{~L}=22,4 \cdot 10^{-3} \mathrm{~m}^{3}$
In $\mathrm{n}=1$ mole of air there is: $\mathrm{V}=22.4 \mathrm{~L}$ (liters)

Table 2. Determination of the universal gas constant - R

| No. crt. | p $\frac{N}{m^{2}}$ | $\begin{aligned} & \mathrm{V} \\ & \mathrm{~m}^{3} \\ & \mathbf{1 0}^{-6} \end{aligned}$ | 年 | $\begin{aligned} & \stackrel{?}{e} \\ & = \\ & \stackrel{\theta}{\theta} \\ & \underset{a}{2} \end{aligned}$ | $\begin{aligned} & R=\frac{p \cdot V}{n \cdot T} \\ & \frac{J}{m o l . K} \end{aligned}$ | $\begin{aligned} & \bar{R} \\ & \frac{J}{\mathrm{~mol} . K} \end{aligned}$ | $\mu= \pm \sqrt{\frac{\sum\left(R_{i}-\bar{R}\right)^{2}}{n(n-1)}}$ | Rreal |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 1 |  |  |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |  |
| 4 |  |  |  |  |  |  |  |  |
| 5 |  |  |  |  |  |  |  |  |

Average values, errors and interval of real values can be calculated:
Average value for $\overline{\mathrm{R}}=\frac{\sum_{i=1}^{n} R_{i}}{n}$
The mean squared error of the average: $\mu= \pm \sqrt{\frac{\sum\left(R_{i}-\bar{R}\right)^{2}}{n(n-1)}}$

## DETERMINATION OF SPECIFIC HEAT FOR SOLID AND LIQUID MATERIALS

## INTRODUCTION

In order to determine the temperature of different materials we should know the amount of heat given off or absorbed by each material within the system. This depends on the mass of the material, the nature of the material from which the body is made and the degree of heating. The heating of a material can be achieved by radiation, convection currents and thermal conductivity.

The thermal capacity of a material is equal to the heat required by that material to increase or decrease its temperature by one degree Kelvin $C=\frac{Q}{T}$

The unit of measurement is: $[C]_{\text {S.I. }}=\frac{[Q]_{\text {S.I. }}}{[T]_{\text {S.I. }}}=\frac{\mathrm{J}}{\mathrm{K}}$ or $[C]_{\text {S.I. }}=\frac{\mathrm{kcal}}{\mathrm{K}}$
THERMAL CAPACITY of the mass unit is called specific heat.
$c=\frac{C}{m}$ whose unit of measurement is: $[c]_{\text {S.I. }}=\frac{\mathrm{J}}{\mathrm{kgK}}$ or $[c]_{\text {S.I. }}=\frac{\mathrm{kcal}}{\mathrm{kgK}}$
SPECIFIC HEAT is the heat required by the mass unit of the material increase or decrease its temperature by 1 degree $c=\frac{Q}{m \cdot \Delta T}$

And the THERMAL CAPACITY of a mole of substance is called molar heat and is the amount of heat needed by a mole of substance to increase or decrease its temperature by 1 degree Kelvin.
$c_{\mu}=\frac{Q}{v \cdot \Delta T}$ the unit of measurement is: $\left[c_{\mu}\right]_{\text {S.I. }}=\frac{J}{\mathrm{kmolK}}$ or $\left[c_{\mu}\right]_{\text {S.I. }}=\frac{\mathrm{kcal}}{\mathrm{kmolK}}$
Between the international system of units and the tolerated system the following relationship exists: $1 \frac{\mathrm{kcal}}{\mathrm{K}}=4.186 \frac{\mathrm{~J}}{\mathrm{~K}}$

The specific heat is a variable quantity with the temperature generally increasing with its raise, therefore its definition must specify the temperature to which it refers. For solid bodies heated in the temperature range between $273 K$ and $373 K$ the specific heat is measured by the mixture method
based on the thermal balance with the help of the calorimeter. The mixture method is based on two calorimetric principles:

## 1. The principle of heat exchange equality:

When there is an exchange of heat and only heat between the two systems of bodies, until thermal equilibrium is established, the sum of the heat received by the system initially at the lower temperature is equal to the sum of the heat lost by the system of bodies initially at a higher temperature raised.

## 2. Principle of inverse transformations:

The amount of heat required for a system of bodies to go from a state 1 to a state 2 is that given by the system passing from state 2 to state 1 . The method of mixtures consists in heating the considered body to a known temperature and introducing it inside the vessel with water, measuring the increase in its temperature. If it is noted with $C$ the thermal capacity of the calorimeter, with $M_{a}$ and $c_{a}$ - the mass, respectively the specific heat of the water in the calorimeter, $T_{0}$ - its initial temperature, $T_{f}$ - the final temperature of the mixture.

After the body to be studied has been introduced into the calorimeter considering that: $Q=M c \Delta T$ (1)
the calorimetric equation is written: $\left(c_{a} M_{a}+C\right)\left(T_{f}-T_{0}\right)=M c\left(T-T_{f}\right)$
where the specific heat of the body to be studied is obtained: $c=\frac{\left(C+M_{a} c_{a}\right)\left(T_{f}-T_{0}\right)}{M\left(T-T_{f}\right)}$
The thermal capacity of the calorimeter is given by the relation: $C=\sum_{i=1}^{n} M_{i} c_{i}$
$M_{i}, c_{i}$ - represent the masses and the specific heats respectively and the indexes have the following specification: 1 - refers to the calorimetric vessel, 2 - stirrer, 3 - thermometer.

The purpose of this paper is to determine the specific heat for a metallic body using the Berthelot calorimeter.

## EXPERIMENTAL ASSEMBLY:


calorimetru

The necessary equipment consists of a Berthelot calorimeter, an electric hob and a vessel for heating the bodies that will be studied. The Berthelot calorimeter has the following components: isolating cover, thermometer, stirrer, calorimetric vessel.

Figure 1. Schematic representation of the experimental assembly for determination of specific heat

## DETERMINATION OF HEAT CAPACITY OF THE CALORIMETER

## PROCEDURE FOR THE PRACTICAL WORK - STEPS AND EXPERIMENTAL DATA PROCESSING

## The calorific capacity of the calorimeter is determined as follows:

- A known amount of water $\mathrm{Ma}_{1}$ is introduced in the calorimeter, the temperature of the water should be in thermal equilibrium with the ambient temperature - noted as temperature $\mathrm{t}_{0}$
- Afterwards a known amount of heated water $\mathrm{Ma}_{2}$ is added in the calorimeter, noted as temperature t
- The final temperature of the mixture $\left(\mathrm{t}_{\mathrm{e}}\right)$ is determined by using a thermometer
- The obtained data will be inserted in table 1.
- Calculations are made for the heat capacity of the calorimeter

$\mathrm{M}_{\mathrm{a} 1}$ - mass of cold water introduced in calorimeter
$\mathrm{M}_{\mathrm{a} 2}$ - mass of hot water introduced in calorimeter
$c_{a}=4182.61 \frac{\mathrm{~J}}{\mathrm{kgK}}$ - specific heat of water
The calorimetric equation is as follows: $\left|Q_{\text {ced }}\right|=Q_{a b s}$
$Q_{a b s}=M_{a 1} c_{a}\left(t_{e}-t_{0}\right)+C\left(t_{e}-t_{0}\right)>0$
$Q_{c e d}=M_{a 2} c_{a}\left(t_{e}-t\right)<0$
$\left|Q_{c e d}\right|=M_{a 2} c_{a}\left(t-t_{e}\right)>0$
$M_{a 1} c_{a}\left(t_{e}-t_{0}\right)+C\left(t_{e}-t_{0}\right)=M_{a 2} c_{a}\left(t-t_{e}\right)$

$$
\begin{equation*}
C=\frac{M_{a 2} c_{a}\left(t-t_{e}\right)}{\left(t_{e}-t_{0}\right)}-M_{a 1} c_{a} \tag{5}
\end{equation*}
$$

Table 1: Determination of the heat capacity of the calorimeter

| n | $\begin{aligned} & M_{a 1} \\ & (k g) \end{aligned}$ | $\begin{aligned} & M_{a 2} \\ & (k g) \end{aligned}$ | $\begin{aligned} & c_{a} \\ & \left(\frac{J}{k g K}\right) \end{aligned}$ | $\begin{aligned} & t_{0} \\ & \left({ }^{0} \mathrm{C}\right) \end{aligned}$ | $\left({ }^{0} C\right)$ | $\begin{aligned} & \hline t_{e} \\ & \left({ }^{0} C\right) \end{aligned}$ | $\begin{aligned} & C \\ & \left(\frac{J}{K}\right) \end{aligned}$ | $\begin{aligned} & \bar{C} \\ & \left(\frac{J}{K}\right) \end{aligned}$ | $\left(\frac{J}{K}\right)$ | $\begin{gathered} C_{\text {real }} \\ \left(\frac{J}{K}\right) \end{gathered}$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 1 |  |  |  |  |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |  |  |  |
| 3 |  |  | 4182.61 |  |  |  |  |  |  |  |
| 4 |  |  |  |  |  |  |  |  |  |  |
| 5 |  |  |  |  |  |  |  |  |  |  |

Average values, errors and the interval of actual values will be calculated
The average value of the caloric capacity is $\overline{\mathrm{C}}=\frac{\sum_{i=1}^{n} C_{i}}{n}$
The mean squared error of the average: $\mu= \pm \sqrt{\frac{\sum\left(C_{i}-\bar{C}\right)^{2}}{n(n-1)}}$
The interval of real values is: $C_{\text {real }} \in(\bar{C}-\mu, \bar{C}+\mu)$

## DETERMINATION OF THE SPECIFIC HEAT FOR SOLIDS

## The specific heat for solid materials is determined as follows:

- A known quantity of water $\mathrm{M}_{\mathrm{a}}$ that is in thermal equilibrium with the ambient temperature, at temperature $t_{0}$, is introduced into the calorimeter
- After that a solid material of known mass-m (whose mass is determined with a balance before), is introduced in the calorimeter at a specific temperature ( t )
- The final temperature $\left(\mathrm{t}_{\mathrm{e}}\right)$ of the mixture (water-solid material) is determined
- The obtained data will be entered in table 1.
- Calculations are made for the specific heat of the solid material that will be studied.

$M_{a 1}+M_{a 2}=M_{a}$ is the mass of cold water in the calorimeter
$\mathbf{m}$ - is the mass of the solid material whose specific heat (c) will be determined. This solid material (a body made from teflon) is heated to a specific temperature ( $\mathbf{t}$ )

The calorimetric equation is: $\left|Q_{c e d}\right|=Q_{a b s}$
$Q_{a b s}=M_{a} c_{a}\left(t_{e}-t_{0}\right)+C\left(t_{e}-t_{0}\right)>0$
$Q_{\text {ced }}=m . c .\left(t_{e}-t\right)<0$
$\left|Q_{\text {ced }}\right|=m . c .\left(t-t_{e}\right)>0$
$M_{a} c_{a}\left(t_{e}-t_{0}\right)+C\left(t_{e}-t_{0}\right)=m . c\left(t-t_{e}\right)$
$. c=\frac{M_{a} c_{a}\left(t_{e}-t_{0}\right)+C\left(t_{e}-t_{0}\right)}{m \cdot\left(t-t_{e}\right)}$

Table 2: Determination of the specific heat for solid materials

| n | $M_{a}$ $(k g)$ | $\begin{aligned} & m \\ & (k g) \end{aligned}$ | $\int_{\left(\frac{J}{k g K}\right)}^{c_{a}}$ | $\begin{aligned} & t_{0} \\ & \left({ }^{\circ} \mathrm{C}\right) \end{aligned}$ | $\begin{aligned} & t \\ & \left({ }^{0} C\right) \end{aligned}$ | $\begin{array}{\|l\|} \hline t_{e} \\ \left({ }^{\circ} C\right) \end{array}$ | $\begin{aligned} & \hline c \\ & \left(\frac{J}{K g \cdot K}\right) \end{aligned}$ | $\begin{aligned} & \bar{c} \\ & \left(\frac{J}{K g . K}\right) \end{aligned}$ | $\left(\frac{J}{K g . K}\right)$ | $\begin{gathered} c_{\text {real }} \\ \left(\frac{J}{\mathrm{Kg} \cdot \mathrm{~K}}\right) \end{gathered}$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 1 |  |  |  |  |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |  |  |  |
| 3 |  |  | 4182.61 |  |  |  |  |  |  |  |
| 4 |  |  |  |  |  |  |  |  |  |  |
| 5 |  |  |  |  |  |  |  |  |  |  |

Average values, errors and the interval of actual values will be calculated
Average values of $\mathrm{c}: \overline{\mathrm{c}}=\frac{\sum_{i=1}^{n} c_{i}}{n}$ and
The mean squared error of the average: $\mu= \pm \sqrt{\frac{\sum\left(c_{i}-\bar{c}\right)^{2}}{n(n-1)}}$
The real value will be in the following interval : $c_{\text {real }} \in(\bar{c}-\mu, \bar{c}+\mu)$

## DETERMINATION OF THE SPECIFIC HEAT OF LIQUIDS

## The specific heat of liquids is determined as follows:

- A known amount of lichd $\mathrm{M}_{1}$ is introduced into the calorimeter in thermal equilibrium with the ambient temperature, at temperature $\mathrm{t}_{0}$
- Afterwards a known quantity of the same heated liquid $\left(\mathrm{M}_{2}\right)$ at temperature t is added into the calorimeter
- The final temperature of the mixture $t$ is determined.
- The obtained data will be entered in table 1 .
- Calculations are made for the specific heat of the studied liquid $c_{L}$

$\mathrm{M}_{1}$ - mass of cold liquid in the calorimeter
$\mathrm{M}_{2}$ - mass of hot liquid in the calorimeter
$\mathrm{c}_{\mathrm{L}}$ - specific heat of the liquid that will be studied
The calorimetric equation is: $\left|Q_{c e d}\right|=Q_{a b s}$
$Q_{a b s}=M_{1} c_{L}\left(t_{e}-t_{0}\right)+C\left(t_{e}-t_{0}\right)>0$
$Q_{c e d}=M_{2} c_{L}\left(t_{e}-t\right)<0$
$\left|Q_{\text {ced }}\right|=M_{2} c_{L}\left(t-t_{e}\right)>0$
$M_{1} c_{L}\left(t_{e}-t_{0}\right)+C\left(t_{e}-t_{0}\right)=M_{2} c_{L}\left(t-t_{e}\right)$
$. c_{L}=\frac{C\left(t_{e}-t_{0}\right)}{M_{2} \cdot\left(t-t_{e}\right)-M_{2}\left(t_{e}-t_{0}\right)}$

Table 3: Determination of the specific heat of a liquid

| n | $\begin{aligned} & M_{1} \\ & (\mathrm{~kg}) \end{aligned}$ | $\begin{aligned} & M_{2} \\ & (\mathrm{~kg}) \end{aligned}$ | C <br> (J/K) | $t_{0}$ $\left({ }^{0} C\right)$ | $\left({ }^{t} C\right)$ | $\begin{aligned} & t_{e} \\ & \left({ }^{0} C\right) \end{aligned}$ | $\begin{aligned} & c_{L} \\ & \left(\frac{J}{K g . K}\right) \end{aligned}$ | $\begin{aligned} & \bar{c}_{L} \\ & \left(\frac{J}{K g . K}\right) \end{aligned}$ | $\left(\frac{J}{K g \cdot K}\right)$ | $\begin{gathered} c_{\text {L real }} \\ \left(\frac{J}{K g \cdot K}\right) \end{gathered}$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 1 |  |  |  |  |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |  |  |  |
| 4 |  |  |  |  |  |  |  |  |  |  |
| 5 |  |  |  |  |  |  |  |  |  |  |

Average values, errors and the interval of actual values will be calculated
Average value of $\mathrm{c}_{\mathrm{L}}: \overline{\mathrm{c}}_{\mathrm{L}}=\frac{\sum_{i=1}^{n} c_{i}}{n}$ and
The mean squared error of the average: $\mu= \pm \sqrt{\frac{\sum\left(c_{i}-\bar{c}\right)^{2}}{n(n-1)}}$
The real value will be in the interval : $c_{L \text { real }} \in\left(\bar{c}_{L}-\mu, \bar{c}_{L}+\mu\right)$

## THE STUDY OF FREE DIFFUSION

## INTRODUCTION

Dilute solutions present similar characteristics with gas mixtures.
The molecules of the dissolved substance diffuse, meaning that, they spread throughout the mass of the solvent liquid.

By dissolving a solid in a liquid, the ordered structure is lost, and its molecules get a supplement of kinetic energy, being hit by the molecules of the solvent.

In turn, the solvent molecules lose part of their kinetic energy, which is expressed by the decrease in the temperature of the solution. In order to restore the initial temperature, heat must be supplied.

Diffusion transports molecules (or ions) from a higher concentration to a lower concentration.

The spontaneous variation in time of the concentrations of the components of a system due to the relative movement of their particles, consists of a transport of mass or electric charge.

Diffusion in gases or in liquids consists in the spreading of the molecules of a gas (liquid) among the molecules of another gas (liquid), when the substances come into contact, resulting in a mass transport in the sense of decreasing the concentration of the molecules of each of the considered components.

According to $\underline{\text { Fick's law: }} \quad \frac{d n}{d t \cdot S}=-D \frac{d C}{d x}$ or $\frac{\Delta n}{\Delta t \cdot S}=-D \frac{\Delta C}{\Delta x}$
Where J is the material flow

- On the one hand, the flow of material is equal to the concentration gradient $\frac{d C}{d x}$ or concentration variation $\Delta \mathrm{c}$ in a given direction x , i.e. $\Delta \mathrm{x}$

$$
J=-D \frac{d C}{d x} \text { or } J=-D \frac{\Delta C}{\Delta x}
$$

- However, the Material Flow is the number of moles of solute $\Delta \mathrm{n}$ transported due to diffusion through an area $S$ of $1 \mathrm{~m}^{2}$ in a time interval of $\Delta \mathrm{t}=1 \mathrm{~s}$.
$J=\frac{\Delta n}{\Delta t . S}$
$[J]_{\text {S.I. }}=1 \frac{\mathrm{~mol}}{\mathrm{~m}^{2} \mathrm{~s}}$
$D$ is the diffusion coefficient which is a characteristic of the dissolved substance but also of the environment in which the diffusion takes place. The scientific literature provides tabulated data for the diffusion coefficient $D$.

Taking into account the formula: $J=-D \frac{\Delta C}{\Delta x}$, the unit of measurement for the concentration gradient
$\frac{d C}{d x}$ is: $\left[\frac{\Delta C}{\Delta x}\right]_{\text {S.I. }}=1 \frac{\mathrm{~mol}}{\mathrm{~m}^{3} \cdot \mathrm{~m}}$
where $\frac{m o l}{m^{3}}$ is the unit of measure for the concentration
Starting from the equation $J=-D \frac{\Delta C}{\Delta x}$, from the point of view of the unit of measurement, the following unit was obtained:

$$
\frac{\mathrm{mol}}{\mathrm{~m}^{2} . S .}=[D]_{\text {S.I. }} \cdot \frac{\mathrm{mol}}{\mathrm{~m}^{4}} \Rightarrow[D]_{\text {S.I. }}=1 \frac{\mathrm{~m}^{2}}{\mathrm{~s}}
$$

Fick's law is valid only for stationary conditions! (meaning that, for boundary conditions in which it can be considered that a very small value of the flow - in a very small time interval - does not significantly change the concentration gradient).
Let's consider 2 vessels of infinite size that communicate with each other through a channel, through which there is no tendency for the solute to flow from the region with higher concentration to the one with lower concentration. If the liquid level in the two compartments are equal then, from the point of view of the law of communicating vessels, there is no mechanical flow tendency from one compartment to another.
To measure the concentration variation in the studied solution, we measure the electrical conductivity. Electrical conductivity is the inverse of resistance: $G=\frac{1}{R}$

The unit of measurement for conductivity is $[G]_{\text {S.I. }}=\frac{1}{\Omega}=S($ SIEMENS $)$


For a resistor of length 1 , section $S$ and resistivity $\rho$, the expression
 of electrical resistance can be written: $R=\frac{\rho \cdot l}{S}$.
The unit of measurement for resistivity is: $[\rho]_{S . I .}=1 \Omega . \mathrm{m}$.
We can also talk about specific conductivity: $\lambda=\frac{1}{\rho}$

## EXPERIMENTAL PART

An aqueous solution of KCl is used.
Electrolytes can be classified as "strong" or "weak". Strong electrolytes dissociate completely into ions upon dissolution, while weak electrolytes dissociate only partially upon dissolution.
$\mathbf{K C l}$ is an electrolyte that dissociates as follows: $\mathrm{KCl} \rightarrow \mathrm{K}^{+}+\mathrm{Cl}^{-}$
Electrolytes are ionic forms that dissolve in water
All salts are weak electrolytes: $\mathrm{KCl}, \mathrm{NaCl}, \mathrm{CH}_{3} \mathrm{COONa}$
Non-electrolytes are substances that dissolve in water and do not dissociate, such as: glucose $\mathrm{C}_{6} \mathrm{H}_{12} \mathrm{O}_{6}$, urea - $\mathrm{NH}_{2} \mathrm{CONH}_{2}$

Examples of strong and weak electrolytes

| ACID |  | BASES |  |
| :---: | :---: | :---: | :---: |
| Strong acids | Weak acids | Strong bases | Weak bases |
| SOLID INORGANIC ACIDS nitric acid, $\mathrm{HNO}_{3}$ hydrobromic acid, HBr hydrochloric acid, HCl hydrochloric acid, $\mathrm{HClO}_{3}$ hydroiodic acid, HI perchloric acid, $\mathrm{HClO}_{4}$ sulfuric acid, $\mathrm{H}_{2} \mathrm{SO}_{4}$ | WEAK AND MEDIUM STRENGTH INORGANIC ACIDS <br> boric acid, $\mathrm{H}_{3} \mathrm{BO}_{3}$ <br> carbonic acid, $\mathrm{H}_{2} \mathrm{CO}_{3}$ (simply called "acid" when included in mineral waters or soft drinks such as Fanta) <br> hydrofluoric acid, HF <br> WEAK ORGANIC ACIDS <br> acetic acid, $\mathrm{CH}_{3} \mathrm{COOH}$ <br> butyric acid, $\mathrm{CH}_{3} \mathrm{CH}_{2} \mathrm{CH}_{2} \mathrm{COOH}$ <br> benzoic acid, $\mathrm{C}_{6} \mathrm{H}_{5} \mathrm{COOH}$ <br> citric acid, $\left(\mathrm{C}_{6} \mathrm{O}_{7} \mathrm{H}_{8}\right)$ <br> formic acid, HCOOH <br> lactic acid, $\mathrm{CH}_{3} \mathrm{CHOHCOOH}$ <br> malic acid, $\mathrm{HOOCCH}_{2} \mathrm{CHOHCOOH}$ <br> oxalic acid, HOOCCOOH <br> pyruvic acid, $\mathrm{CH}_{3} \mathrm{COCOOH}$ <br> propionic acid, $\mathrm{CH}_{3} \mathrm{CH}_{2} \mathrm{COOH}$ <br> valerianic acid, $\mathrm{CH}_{3}\left(\mathrm{CH}_{2}\right)_{3} \mathrm{COOH}$ | NaOH <br> KOH <br> LiOH <br> $\mathrm{Ca}(\mathrm{OH})_{2}$ $\mathrm{Ba}(\mathrm{OH})_{2}$ | $\begin{aligned} & \mathrm{Ammonia} \mathrm{NH}_{3} \\ & \mathrm{H}_{2} \mathrm{O} \\ & \mathrm{Mg}(\mathrm{OH})_{2} \\ & \mathrm{Zn}(\mathrm{OH})_{2} \\ & \mathrm{Al}(\mathrm{OH})_{3} \end{aligned}$ |

DESCRIPTION OF THE WORKING INSTRUMENT: MULTIMETER WTW 340I (CONDUCTOMETER / PH METER):


1- the keyboard
2- display
3 - current connection
Affichage

|  | affichage de l'état |
| :---: | :---: |
| $\mathrm{UPHO}_{2} \mathrm{XS}^{3} \mathrm{al}{ }^{\text {a }}$ | symbole de senseur / l'électrode |
|  |  |
|  | affichage de la valeur mesurée |
|  | fonction et affichage de la température |
|  |  |


|  | M | sélectionner la variable mesurée pH / potentiel Redox D.O. concentration / D.O. saturation conductivité / salinité |
| :---: | :---: | :---: |
|  | sто | enregistrer (sauver) une valeur mesurée |
|  | $\checkmark$ | ouvert / fermé ( ON/OFF) |
| $\nabla$ | CAL | calibrer la variable mesurée actuellement fixé |
|  | RCL | Afficher / transmettre les valeurs mesurées |
|  | AR | activer / désactiver la fonction de lecture automatique (AutoRead) |
|  | $\Delta$ | sélectionner le mode de mesure, augmenter les valeurs, rouleau |
|  | $\nabla$ | sélectionner le mode de mesure, diminuer les valeurs, rouleau |
|  | RUNENTER | confirmer un enregistrement, activer la fonction lecture automatique (AutoRead) valeurs de sortie mesurée |




1. Put distilled water (AD) in glass 3

Distilled water dissociates as follows: $2 \mathrm{H}_{2} \mathrm{O} \Leftrightarrow \mathrm{H}_{3} \mathrm{O}^{+}+\mathrm{OH}^{-}$
Where $\mathrm{H}_{3} \mathrm{O}^{+}$- is called the hydronium ion and $\mathrm{OH}^{-}$the hydroxyl ion.
Water does not dissociate well because it is not a strong electrolyte
2. Within the space created (2) by the two conductivity probes, a porous material is inserted - for example cotton (wadding)
3. The conductivity of the distilled water is measured

The conductance is read $G\left(\frac{\mu S}{c m}\right)$ every 30 seconds until 5 minutes when the value of G stabilizes
4. The recorded data will be entered in table 1 .
5. Change the aqueous solution: distilled water from glass (3) with $1 \% \mathrm{KCl}$ aqueous solution.
6. Measurements are made every 30 seconds up to 30 minutes.
7. Do not stir the solution to observe only the phenomenon of diffusion

Table 1

| t (s) | $\mathrm{G}(\mu \mathrm{S} / \mathrm{cm})$ | $\mathrm{c}\left(\mathrm{mol} / \mathrm{m}^{3}\right)$ | G/c |
| :---: | :---: | :---: | :---: |
| 0 |  |  |  |
| 30 |  |  |  |
| 60 |  |  |  |
| 90 |  |  |  |
| 120 |  |  |  |
| 150 |  |  |  |
| 180 |  |  |  |
| 210 |  |  |  |
| 240 |  |  |  |
| 270 |  |  |  |
| 300 |  |  |  |
| 330 |  |  |  |
| 360 |  |  |  |
| 390 |  |  |  |
| 420 |  |  |  |
| 450 |  |  |  |
| 480 |  |  |  |
| 510 |  |  |  |
| 540 |  |  |  |
| 570 |  |  |  |
| 600 |  |  |  |
| 630 |  |  |  |
| 660 |  |  |  |
| 690 |  |  |  |
| 720 |  |  |  |
| 750 |  |  |  |
| 780 |  |  |  |
| 810 |  |  |  |
| 840 |  |  |  |
| 870 |  |  |  |
| 900 |  |  |  |
| 930 |  |  |  |
| 960 |  |  |  |
| 990 |  |  |  |


| 1020 |  |  |  |
| :--- | :--- | :--- | :--- |
| 1050 |  |  |  |
| 1080 |  |  |  |
| 1110 |  |  |  |
| 1140 |  |  |  |
| 1170 |  |  |  |
| 1200 |  |  |  |
| 1230 |  |  |  |
| 1260 |  |  |  |
| 1290 |  |  |  |
| 1320 |  |  |  |
| 1350 |  |  |  |
| 1380 |  |  |  |
| 1410 |  |  |  |
| 1440 |  |  |  |
| 1470 |  |  |  |
| 1500 |  |  |  |
| 1530 |  |  |  |
| 1560 |  |  |  |
| 1590 |  |  |  |
| 1620 |  |  |  |
| 1650 |  |  |  |
| 1680 |  |  |  |
| 1710 |  |  |  |
| 1740 |  |  |  |
| 1770 |  |  |  |
| 1800 |  |  |  |
|  |  |  |  |

8. Open an Excel file and fill in table 1, enter the conductance G measured at every 30 s .
9. draw the graph $G=G(t)$. It is a graph that after a period of instability of approx. 3 minutes shows a linear variation of conductivity depending on time. The linear part of the curve is selected and the equation corresponding to this part is determined: $\mathbf{G}=\mathbf{A}_{\mathbf{1}}^{\mathbf{1}} \mathbf{t}+\mathbf{B} \mathbf{1}$. The value of the coefficient $\mathrm{R}^{2}$ should be close to the value 1 . If the value of $\mathrm{R}^{2}$ is closest to 1 , this means that this is the line that best represents the phenomenon.
$\mathrm{A}_{1}$ is the slope of the equation line: $\mathbf{G}=\mathbf{A}_{\mathbf{1}} . \mathbf{t}+\mathbf{B}_{\mathbf{1}}$.

So we represented the variation of the conductance as a function of time: $\mathrm{G}=\mathrm{G}(\mathrm{t})$
10. To determine the variation in the concentration of ions recorded in the region where the probe was inserted, do not remove the porous layer of cotton between the two conductometric sensors and measure the conductivity of a $0.1 \% \mathrm{KCl}$ solution. The correlation between conductivity and concentration is made.

For the calibration solution ( $\mathrm{KCl} 0.1 \%$ ) the following value is used $G=1923 \frac{\mu \mathrm{~S}}{\mathrm{~cm}}$
The molar mass of KCl is: $M_{K C l}=74,5 \frac{\mathrm{~g}}{\mathrm{~mol}}$
For the concentration of $0.1 \% \mathrm{KCl}$ we have the measured value of $G=1923 \frac{\mu \mathrm{~S}}{\mathrm{~cm}}$
For the concentration of $1000 \frac{\mathrm{~g}}{\mathrm{~m}^{3}} \mathrm{KCl}$ the measured value of G is as follows $G=1923 \frac{\mu \mathrm{~S}}{\mathrm{~cm}}$
For the concentration of $\frac{1000}{74.5} \frac{\mathrm{~mol}}{\mathrm{~m}^{3}} \mathrm{KCl}$ the measured value of G is as follows $G=1923 \frac{\mu \mathrm{~S}}{\mathrm{~cm}}$
Then for the unknown concentration $c\left(\frac{\mathrm{~mol}}{\mathrm{~m}^{3}}\right) K C l$ we have the measured value of $G\left(\frac{\mu \mathrm{~S}}{\mathrm{~cm}}\right)$
$\operatorname{so} c\left(\frac{\mathrm{~mol}}{\mathrm{~m}^{3}}\right)=G \cdot \frac{1000 \frac{\mathrm{~mol}}{\mathrm{~m}^{3}}}{74,5 \frac{\mathrm{~g}}{\mathrm{~mol}} \cdot 1923 \frac{\mathrm{\mu S}}{\mathrm{~cm}}}=G \cdot 0,006980145$
$c=G \cdot \frac{1000}{74,5 \cdot 1923}=G \cdot 0,006980145$
This represents the formula to calculate the concentration
For the values obtained in this formula, enter the values obtained in the second column of table 1.
11. Draw the graph $\mathrm{c}=\mathrm{c}(\mathrm{t})$ as a function of time.

The line is obtained: $\mathrm{c}=\mathrm{A}_{2} \cdot \mathrm{t}+\mathrm{B}_{2}$ with the $\mathrm{R}^{2}$ coefficient also very close to 1 .
12. The flow of ions is calculated using Fick's law: $J=\frac{\Delta n}{\Delta t \cdot S}=\frac{\Delta n \cdot v}{v \Delta t \cdot S}$.
but $\frac{\Delta n .}{v}=\Delta C$ - it is the variation of the concentration, that is, the variation of the number of moles in the volume unit.
then the particle flow can be written: $J=\frac{\Delta n}{\Delta t \cdot S}=\frac{\Delta c}{\Delta t} \cdot \frac{v}{S}$
$\frac{\Delta c .}{\Delta t}=A_{2}$ represents the slope of the line $\mathrm{c}=\mathrm{A}_{2} \cdot \mathrm{t}+\mathrm{B}_{2}$ which is the variation concentration as a function of time $\mathrm{c}=\mathrm{c}(\mathrm{t})$.

In the flow expression $J=\frac{\Delta n}{\Delta t \cdot S}=\frac{\Delta c \cdot}{\Delta t} \cdot \frac{v}{S}$, the ratio between volume and surface is: $\frac{v}{S}=\Delta x=1 \mathrm{~cm}=10^{-2} \mathrm{~m}$.

It is considered $S=1 \mathrm{~cm}^{2}$ the surface of the conductor which is $S=1 \mathrm{~cm}^{2}$ is subject to the diffusion phenomenon and the corresponding volume $\mathrm{V}=1 \mathrm{~cm}^{3}$
Remember: the volume of the probe is used here, which is $V=1 \mathrm{~cm}^{3}$

Thus the flux corresponding to the porous material is $\mathbf{J}=\mathbf{A}_{\mathbf{2} .10}{ }^{-\mathbf{2}}$

It is worth noting here the significance of the notation $\mathrm{A}_{2}$
$\frac{\Delta c .}{\Delta t}=A_{2}$ represents the slope of the line $\mathrm{c}=\mathrm{A}_{2} \cdot \mathrm{t}+\mathrm{B}_{2}$ which represents concentration as a function of time $\mathrm{c}=\mathrm{c}(\mathrm{t})$

According to Fick's equation, the ion flux that diffuses is $J=-D \cdot \frac{\Delta c}{\Delta x}$
But here we have a concentration difference from $1 \%$ to $0 \%$.
Then: $\cdot \frac{\Delta c}{\Delta x}=134,2 \frac{\mathrm{~mol}}{\mathrm{~m}^{3}} \cdot \frac{1}{10^{-2} \mathrm{~m}}=134,2 \cdot 10^{2} \frac{\mathrm{~mol}}{\mathrm{~m}^{4}}$
It can be said that the concentration equivalent of $1 \%$ corresponds to a value of $134,2 \frac{\mathrm{~mol}}{\mathrm{~m}^{3}}$
$1 \%---\frac{1 \mathrm{~g}}{100 \mathrm{ml}}---10 \frac{\mathrm{~g}}{\mathrm{l}}---10000 \frac{\mathrm{~g}}{\mathrm{~m}^{3}}---\frac{10000}{74,5} \frac{\mathrm{~mol}}{\mathrm{~m}^{3}}---134,2 \frac{\mathrm{~mol}}{\mathrm{~m}^{3}}$
It can be written as follows: $\mathrm{A}_{2}$ corresponds to the slope of the concentration line $\mathrm{c}=\mathrm{c}(\mathrm{t})$
$A_{2} \cdot 10^{-2}=-D .134,2 \frac{\mathrm{~mol}}{\mathrm{~m}^{3}}$
$D=\frac{A_{2} \cdot 10^{-2}}{134,2 \cdot 10^{2}}=A_{2} \frac{10^{4}}{134,2} \frac{m^{2}}{s}$ - THE FORMULA FOR CALCULATING THE DIFFUSION

## COEFFICIENT

Where $\mathrm{A}_{2}$ is the slope of the line obtained from the graphic representation

$$
\mathrm{c}=\mathrm{c}(\mathrm{t})
$$

The experiment can be realized with rehydration salts with the composition: NaKCl
The molar mass of NaKCl is: $M_{K N a C l}=97.5 \frac{\mathrm{~g}}{\mathrm{~mol}}$
For rehydration salts, the concentration is calculated with the formula:
$c=G \cdot \frac{1000}{97.5 .1923}$
To determine the diffusion coefficient $D$
The concentration equivalent of $1 \%$ corresponds to a value of


It can be written: $\mathrm{A}_{2}$ corresponds to the slope of the concentration line $\mathrm{c}=\mathrm{c}(\mathrm{t})$
$A_{2} .10^{-2}=-$ D. $102.564 \frac{\mathrm{~mol}}{\mathrm{~m}^{3}}$

$$
D=\frac{A_{2} \cdot 10^{-2}}{102 \cdot 546 \cdot 10^{2}}=A_{2} \cdot \frac{10^{4}}{102.564} \frac{m^{2}}{\mathrm{~s}}
$$

## THE FORMULA FOR CALCULATING THE DIFFUSION COEFFICIENT

Where A2 is the slope of the line obtained from the graphic representation $\mathrm{c}=\mathrm{c}(\mathrm{t})$

## The experiment can be realized with DISTILLED WATER:

The molar mass of $\mathrm{H}_{2} \mathrm{O}$ is: $M_{H 2 O}=18 \frac{\mathrm{~g}}{\mathrm{~mol}}$
For rehydration salts, the concentration is calculated with the formula:
$c=G \cdot \frac{1000}{18.1923}$

## To determine the diffusion coefficient $\mathbf{D}$

The concentration equivalent of $1 \%$ corresponds to a value of


It can be written: $\mathrm{A}_{2}$ corresponds to the slope of the concentration line
$\mathrm{c}=\mathrm{c}(\mathrm{t})$
$A_{2} \cdot 10^{-2}=-$ D. $555,555 \frac{\mathrm{~mol}}{\mathrm{~m}^{3}}$

$$
D=\frac{A_{2} \cdot 10^{-2}}{555.555 \cdot 10^{2}}=A_{2} \cdot \frac{10^{4}}{555.555} \frac{m^{2}}{\mathrm{~s}}
$$

THE FORMULA FOR CALCULATING THE DIFFUSION COEFFICIENT
Where A2 is the slope of the line obtained from the graphic representation $\mathrm{c}=\mathrm{c}(\mathrm{t})$

The rule of determination in the category is the different solutions used, using the Diffusion Coefficient

For MICROMOLECULES:
M<10 ${ }^{\mathbf{3}}$ Da et $\mathrm{D} \sim 10^{-5} \mathrm{~cm}^{2} / \mathrm{s} \sim 10^{-9} \mathrm{~m}^{2} / \mathrm{s}$

## For MACROMOLECULES:

$10^{\mathbf{3}} \mathrm{Da}<\mathrm{M}<10^{8} \mathrm{Da}$ et $\mathrm{D} \sim 10^{-8} \mathrm{~cm}^{2} / \mathrm{s} \mathrm{D} \sim 10^{-12} \mathrm{~m}^{2} / \mathrm{s}$

## STUDY OF VARIATION OF CONDUCTANCE WITH CONCENTRATION FOR STRONG AND WEAK ELECTROLYTES

## THEORETICAL INTRODUCTION

The chemical transformations accompanied by the passage of an electric current through the considered system are known as electrochemical reactions. Many chemical reactions take place by the transfer of electric charges, such as proton transfer reactions between an acid ( AH or $\mathrm{BH}^{+}$) and solvent SH , or between a base ( B or $\mathrm{A}^{-}$) and a solvent S :
$\mathrm{AH}+\mathrm{SH}=\mathrm{A}^{-}+\mathrm{SH}_{2}{ }^{+}$

$$
\mathrm{A}^{-}+\mathrm{SH}=\mathrm{AH}+\mathrm{SB}^{+}
$$

$\mathrm{SH}=\mathrm{S}^{-}+\mathrm{BH}^{+}$

$$
\mathrm{BH}++\mathrm{SH}=\mathrm{B}+\mathrm{SH}_{2}{ }^{+}
$$

In solution, the reactant components are usually in ionic form, but they can also be neutral molecules.

An important role in electrochemical processes is played by the state of ionic species in the liquid phase (for example, in the solution). As electrically charged particles, they interact strongly both with the solvent (an interaction called solvation) and with the other ionic particles (which constitute the ionic atmosphere). Their transition into other forms, following electrochemical reactions, is dependent on the value of the free energy of the process. Both the solvation process (ion - molecule interactions) and the interaction with the ionic atmosphere (ion - ion interactions) present great difficulties in a quantitative approach at the microscopic level.

Electronic conductors are generally solid or molten metals, alloys and some metallic compounds. Electronic conduction occurs through the orderly movement of electrons, while the positive ions from which they came occupy fixed positions. Electronic conduction can be characterized quantitatively with the help of specific conductance, also called conductivity. According to Ohm's law, electrical resistance, $R$, represents the ratio between the applied potential difference, $U$, and the resulting current, $I: R=\frac{\Delta U}{I}$ and is measured in ohms ( $\Omega$ )

The electrical resistance of a body depends on its length, $x$, and its cross-sectional area $S$ (when the section is constant). The proportionality constant, depending on the nature of the material, is called specific resistance or resistivity: $R=\frac{\rho \cdot l}{S}$

The inverse value of resistance is called conductance, and the inverse value of resistivity
$G=\frac{1}{\rho}$ is called conductivity or specific conductance: $\frac{1}{R}=G . \frac{S}{l}$
The dimensions of resistivity and conductivity result from the analysis of the definition relations:
$[\rho]_{\text {S.I. }}=1 \Omega . m \quad ;[G]_{\text {S.II }}=1 \Omega^{-1} \cdot m^{-1}$
The unit $\Omega \quad$ is sometimes referred to as mho (inverse of ohm) or siemens.
Electronic conductors have resistivities between $10^{-8}$ and $10^{-5} \Omega \mathrm{~m}$, compared to insulators that have resistivities $\mathbf{G}>10^{6} \Omega \mathrm{~m}$.

The increase in temperature leads to a decrease in the conductivity of the electronic conductors due to the increase in the agitation of the "electronic gas".

Ionic conductors (also called conductors of the second kind), in the form of electrolyte solutions or melts, are formed either by the electrolytic dissociation of some polar substances (ionic substances such as $\mathrm{HCl}, \mathrm{CH}_{3} \mathrm{COOH}$, etc.) upon dissolution in polar solvents (for example water), either by dissolving solid substances with an ionic structure (ionophores such as $\mathrm{NaCl}, \mathrm{KNO}_{3}$, etc.) or by melting them. Electrolytes are considered strong when they are completely dissociated, or weak when they are only partially dissociated. The division of electrolytes into the two categories has a relative character since even the dissociated species interact with each other. At higher concentrations, "ion pairs" or even higher aggregates are formed.

Unlike electronic conductors, in which the density of charge carriers (electrons) is an intrinsic characteristic of the material, for ionic conductors in the form of solutions the density of charge carriers can be modified by changing the concentration.

In addition, the transport of electric current through an electrolyte solution is provided by both positive and negative ions. The transport of ions in an electrolyte solution occurs either due to the existence of a concentration gradient, or due to the application of a potential difference, or both stresses.

In order to highlight the conductive characteristics of the solutions of the different electrolytes, a standardization of the concentration is first necessary. Conductivity measured at different
concentrations must be converted to a quantity related to a unit concentration. The ratio between conductivity and molar concentration, called molar conductivity, allows the comparison between the conductivities of solutions of different electrolytes at the same concentration. Molar conductivity is defined as: $G_{m}=\frac{1000 \cdot G}{C_{m}}$.
where $C_{m}$ is the molar concentration of the electrolyte, and the factor 1000 converts the molar concentration from $\mathrm{mol} / \mathrm{L}$ to SI units $\left(\mathrm{mol} / \mathrm{m}^{3}\right)$.

Since the conductive properties of an electrolyte solution are dependent on both the concentrations of the ions and their charges, a standardized property also from the point of view of the valence of the ions is the equivalent conductivity, defined as: $G_{e}=\frac{1000 \cdot G}{C_{e}}$.
where $C_{e}$ represents the normal concentration (in gram equivalents per liter of solution).
In the absence of interactions between ions, the equivalent conductivity should be a constant characteristic of each system. It was found experimentally that, for dilute solutions of strong electrolytes, the equivalent conductivity decreases with increasing concentration, according to a law of the form: $G_{C}=G_{0}-A \cdot \sqrt{C}$. By extrapolating the values of $G_{C}$ at zero concentration (when the interionic interactions cancel), the equivalent conductivities at infinite dilution, $\mathrm{G}_{0}$, are obtained. These are intrinsic properties of charge carriers and are tabulated for individual ions in aqueous solutions at a temperature of $25^{\circ} \mathrm{C}$. They are used in the calculation of other transport properties of electrolyte solutions.

The temperature dependence of the conductivity of electrolyte solutions is also different from that of electronic conductors: an increase in temperature leads to an increase in the conductivity of diluted electrolyte solutions according to an empirical law of the form:

$$
G_{t}=G_{25}\left\lfloor 1+\alpha(t-25)-\beta(t-25)^{2}\right\rfloor
$$

The transport of ions in solutions can occur both due to the application of an electric field and due to the existence of a concentration gradient. In the first case, the transport is called migration, and in the second case, diffusion. The diffusion of neutral particles is described by Fick's laws. For ions, these laws must be modified to take into account the electric field gradient induced by a concentration gradient. The two transport mechanisms can occur separately or together.

Electronic and ionic conduction are two extreme mechanisms that can be found together, cooperating in the transport of electric current.

## PROCEDURE FOR THE PRACTICAL WORK - STEPS

An aqueous solution of KCl is used.
Electrolytes can be "strong" or "weak". Strong electrolytes dissociate completely into ions upon dissolution, while weak electrolytes dissociate only partially upon dissolution.
$\mathbf{K C l}$ is an electrolyte that dissociates as follows: $\mathbf{K C l} \rightarrow \mathbf{K}^{+}+\mathrm{Cl}^{-}$
Electrolytes are substances that, in a molten state or in solution, are electrical conductors. The electrical conductivity of electrolyte solutions is very different (whether they are ionic or covalent combinations) and consequently they have been classified according to this criterion into strong electrolytes and weak electrolytes .
Strong electrolytes are the following:

- salts: $\mathrm{NaCl}, \mathrm{KNO}_{3}, \mathrm{Na}_{2} \mathrm{SO}_{4}, \mathrm{~Pb}\left(\mathrm{NO}_{3}\right)_{2}, \mathrm{CaCl}_{2}, \mathrm{CuSO}_{4}, \mathrm{CH}_{3} \mathrm{COONa}$;
- soluble hydroxides: $\mathrm{NaOH}, \mathrm{KOH}, \mathrm{Ba}(\mathrm{OH})_{2}$;
- mineral acids: $\mathrm{HBr}, \mathrm{HCl}, \mathrm{H}_{2} \mathrm{SO}_{4}, \mathrm{HNO}_{3}, \mathrm{H}_{3} \mathrm{PO}_{4}$;

Weak electrolytes are the following:

- organic acids: $\mathrm{CH}_{3} \mathrm{COOH}, \mathrm{C}_{6} \mathrm{H}_{5} \mathrm{COOH}$;
- inorganic acids: $\mathrm{H}_{2} \mathrm{~S}, \mathrm{H}_{2} \mathrm{CO}_{3}, \mathrm{HCN}$;
- bases: $\mathrm{NH}_{3}$, amines.


## DEVICE DESCRIPTION: WTW 340I MULTIMETER

(CONDUCTOMETER / pH METER):


1- keyboard (Keypad)
2- display
3- connection region (Jack field)
domaine de prise


| $\mathbf{1}$ | D.O. sonde ou la conductivité du cellule |
| :--- | :--- |
| 2 | Electrode de pH |
| $\mathbf{3}$ | Sonde de température de pH |
| $\mathbf{4}$ | Branchez - dans l'unité d'alimentation |
| 5 | RS232 interface série |

Affichage




1. Prepare solutions of acetic acid $\mathrm{CH}_{3} \mathrm{COOH}$ of different concentrations as follows:
2. A stock solution of acetic acid of concentration $\mathrm{c}_{1}=0.4 \mathrm{M}$ is made as follows: 12 g of acetic acid in 500 ml of water.
3. From the stock solution of acetic acid with $c_{1}=0.4 \mathrm{M}$, prepare the other solutions needed in the experiment and whose conductance is measured $G\left(\frac{\mu S}{c m}\right)$ as presented in Table 1.

## Table 1

| Soil <br> no | The volume in ml <br> of acetic acid $\mathrm{CH}_{3} \mathrm{COOH}$ of concentration <br> $\mathrm{c}_{1}=0.4 \mathrm{M}$ needed to be poured into the a <br> graduated flask. <br> From this value, add water up to the volume <br> of 50 ml of the graduated flask. | The concentration <br> obtained in molars <br> $(\mathrm{M})$ |
| :--- | :--- | :--- |
| 1 | 45 | 0.36 |
| 2 | 40 | 0.32 |
| 3 | 35 | 0.28 |
| 4 | 30 | 0.24 |
| 5 | 25 | 0.20 |
| 6 | 20 | 0.16 |
| 7 | 15 | 0.12 |
| 8 | 10 | 0.08 |
| 9 | 5 | 0.04 |
| 10 | 2 | 0.016 |
| 11 | 1 | 0.008 |

4. Thus, 12 working solutions are obtained, having concentrations between 0.4 M and 0.008 M , which are introduced into the Berzelius glass and afterwards their conductance will be measured $G\left(\frac{\mu S}{c m}\right)$. The results are shown in table 2 in column 4.
5. Calculate $\mathrm{G} / \mathrm{c}$ and enter the results in table 2 in column 5.
6. A solution of $\mathrm{KCl}-\mathrm{a}$ strong electrolyte - of concentration $\mathrm{c}_{1}=0.4 \mathrm{M}$ is obtained as follows:
dissolve 14.9 g KCl - in 500 ml water
the molar concentration formula is used $c=\frac{m(g)}{M \cdot V(l)}$
where $\mathrm{m}=14.9 \mathrm{~g}$ is the mass of KCl used
$\mathrm{V}=500 \mathrm{ml}=0.5 \mathrm{~L}$ is the volume of water used to obtain the aqueous solution and $\mathrm{M}=74.5 \mathrm{~g} / \mathrm{mol}$ is the molar mass of KCl
7. From the KCl stock solution with $\mathrm{c}_{1}=0.4 \mathrm{M}$, prepare the other solutions needed in the experiment and afterwards their conductance will be measured $G\left(\frac{\mu S}{c m}\right)$ (see table 3)

Table 2:

| No. <br> Sol. | $\mathrm{C}(\mathrm{M})$ | $G\left(\frac{\mu S}{c m}\right)$ <br> $\mathrm{CH}_{3} \mathrm{COOH}$ | $\mathrm{G} / \mathrm{C}$ <br> 3 | $G\left(\frac{\mu S}{c m}\right)$ <br> COOH | $\mathrm{G} / \mathrm{C}$ <br> KCI |
| :--- | :--- | :--- | :--- | :--- | :--- |
| 1 | 0.4 |  |  |  |  |
| 2 | 0.36 |  |  |  |  |
| 3 | 0.32 |  |  |  |  |
| 4 | 0.28 |  |  |  |  |
| 5 | 0.24 |  |  |  |  |
| 6 | 0.20 |  |  |  |  |
| 7 | 0.16 |  |  |  |  |
| 8 | 0.12 |  |  |  |  |
| 9 | 0.08 |  |  |  |  |
| 10 | 0.04 |  |  |  |  |

8. Open an Excel file and fill in the data from table 2:
9. The data will be graphically represented for both $\mathrm{CH}_{3} \mathrm{COOH}$ (weak electrolyte) and KCl (strong electrolyte) as the following function $\mathrm{G}=\mathrm{G}(\mathrm{c})$
10. Afterwards represent graphically $\frac{G}{c}=\frac{G}{c}(c)$ for $\mathrm{CH}_{3} \mathrm{COOH}$ (weak electrolyte) and KCl (strong electrolyte)
11. Compare the graphs obtained for weak and strong electrolytes.

Table 3:

| No. <br> Sol. | Mass of KCl expressed in grams $\mathrm{m}(\mathrm{g})$ <br> required to achieve the concentration | The concentration <br> obtained in molars (M) |
| :--- | :--- | :--- |
| 1 | 14.9 | 0.4 |
| 2 | 13.41 | 0.36 |
| 3 | 11.92 | 0.32 |
| 4 | 10.43 | 0.28 |
| 5 | 8.94 | 0.24 |
| 6 | 7.45 | 0.20 |
| 7 | 5.96 | 0.16 |
| 8 | 4.47 | 0.12 |
| 9 | 2.98 | 0.08 |
| 10 | 1.49 | 0.04 |

## VARIATION OF ELECTRICAL CONDUCTIVITY OF WATER WITH TEMPERATURE

## INTRODUCTION

The electric current passing through a conductive medium encounters a resistance on its part that manifests itself as a difficulty in the passage of the current. To ensure a current intensity ("I") through the conductor, it is necessary to apply a potential difference (electrical voltage, "U") between the ends of the conductive medium. The more a potential difference is needed to ensure a current intensity, the greater is the electrical resistance ("R") of the conductive medium.


Figure 1. Electrical resistance of an environment

Figure 1 shows an imaginary montage for defining the electrical resistance of an environment. The medium of interest is found between planar electrodes "A" and "B" and uniformly fills the parallelepiped $\mathrm{a}, \mathrm{b}, \mathrm{c}, \mathrm{d}, \mathrm{e}, \mathrm{f}, \mathrm{g}, \mathrm{h}$ limited by these electrodes. If the difference of the potential between electrodes "A" and "B" is " $\mathrm{U}_{\mathrm{AB}}$ " (expressed in volts ), and the intensity of the electric current is "I" (expressed in amperes ), then the electrical resistance " $\mathrm{R}_{\mathrm{AB}}$ " of the medium (expressed in ohms ) is calculated according to relation (1).

$$
\mathrm{R}_{\mathrm{AB}}=\frac{\mathrm{U}_{\mathrm{AB}}}{\mathrm{I}}(1)
$$

The electrical resistance of an environment, in addition to the particularity of the environment, depends on its geometric factors according to relation (2).

$$
\mathrm{R}_{\mathrm{AB}}=\rho \cdot \frac{\mathrm{L}}{\mathrm{~S}}(2)
$$

In relation (2), " L " represents the length and " S " is the area of the section of the path traversed by the electric current between electrodes " A " and " B ". The " $\square$ " coefficient is characteristic of the material of the environment; is the resistivity of the environment (in the International System of Units it is expressed in ohms $\cdot$ meter,$\square \cdot \mathrm{m}$ ).

The value of the resistivity of an environment depends on the temperature: in the case of metallic conductors, the resistivity increases with temperature; in the case of semiconductors and electrolyte solutions, the resistivity decreases with temperature. Figure 2 illustrates the temperature dependence of the resistivity of pure water (expressed in the tolerated unit mega-ohm•centimeter, $M \square \cdot \mathrm{~cm}$ ).


Fig. 2 Dependence of pure water resistivity on temperature
The reciprocal value of resistance "R" is expressed in $(\text { ohm })^{-1},(\square)^{-1}$; this unit of measurement is called " Siemens ", " $S$ ".

The reciprocal of resistivity is electrical conductivity, " K ".

$$
\kappa=\frac{1}{\rho}(3)
$$

If for the electrical resistivity " $\rho$ " the tolerated measurement unit ohm $\cdot$ centimeter,$\Omega \cdot \mathrm{cm}$ is accepted, the electrical conductivity "K" is expressed in "Siemens / cm ", " $S / \mathrm{cm}{ }^{\prime \prime}$. Figure 3 represents the conductivity " G " of pure water as a function of temperature.


Fig. 3 Dependence of pure water conductivity on temperature

## PROCEDURE FOR THE PRACTICAL WORK - STEPS

In a 100 ml Berzelius beaker put 50 ml of double-distilled water, a thermometer ("T"), a magnetic stirring bar ("A") and the electrode assembly ("E") of the probe of a conductometer (Fig. 4 ). The glass is placed on a heating plate with the possibility of magnetic stirring. The temperature indicated by the thermometer and the conductivity displayed by the measuring instrument are recorded. Water heating is started and the change in water conductivity is monitored. The temperature vs. value pairs are recorded. conductivity for every $5{ }^{\circ} \mathrm{C}$ temperature change. The monitoring of the process is continued until the temperature of $75^{\circ} \mathrm{C}$.

Stop heating the water and record, every $5{ }^{\circ} \mathrm{C}$, the conductivity vs. temperature during water cooling. The water conductivity is graphically represented as a function of temperature for both heating and cooling.

Compare the graph obtained with the one shown in Figure 4.
How can the observed difference be explained?


Figure 4. Experimental assembly

Table 1: Variation of Conductance with temperature for water

| No. crt | T(oC) | $G(\mu \mathbf{S} / \mathbf{c m})$ |
| :--- | :--- | :--- |
| 1 | 20 |  |
| 2 | 25 |  |
| 3 | 30 |  |
| 4 | 35 |  |
| 5 | 40 |  |
| 6 | 45 |  |
| 7 | 50 |  |
| 8 | 55 |  |
| 9 | 60 |  |
| 10 | 65 |  |
| 11 | 70 |  |
| 12 |  |  |

Table 2: Variation of Conductance with temperature

| No. <br> crt |  | Solution of NaCl 0.9\% | Rehydration salts |
| :--- | :--- | :--- | :--- |
|  | $\mathrm{t}\left({ }^{\circ} \mathrm{C}\right)$ | $\mathrm{G}(\mu \mathrm{S} / \mathrm{cm})$ | $\mathrm{G}(\mu \mathrm{S} / \mathrm{cm})$ |
| 1 | 25 |  |  |
| 2 | 30 |  |  |
| 3 | 31 |  |  |
| 4 | 32 |  |  |
| 5 | 33 |  |  |
| 6 | 34 |  |  |
| 7 | 35 |  |  |
| 8 | 36 |  |  |
| 9 | 37 |  |  |
| 10 | 38 |  |  |
| 11 | 39 |  |  |
| 12 | 40 |  |  |
| 14 | 41 |  |  |
| 15 | 42 |  |  |
| 16 | 43 |  |  |
| 17 | 44 |  |  |
| 18 | 45 |  |  |

Graphs are represented: $\mathrm{G}=\mathrm{G}(\mathrm{t})$ for tap water (table 1 ) and for solution of $\mathrm{NaCl} 0.9 \%$ and rehydration salts (table 2).

Conclusions are formulated.

## DETERMINATION OF TDS

## THEORETICAL INTRODUCTION

TDS - total dissolved solids is a measure of the dissolved content of all combined inorganic and organic substances that are present in a liquid in the form of a molecular, ionized or microgranular suspension (colloidal soil). In general, the definition of such sizes is that the solids must be small enough to pass through micrometer size (normal size or smaller) filters. TDS are normally discussed only for freshwater systems because salinity includes some of the ions that make up the definition of TDS

The main application of TDS is the study of water quality. Although TDS is not generally considered a primary pollutant (ie, it is not considered to be associated with health effects), it is used as an indication of aesthetic characteristics in drinking water and as an aggregate indicator of the presence of a wide range of chemical contaminants.

The primary sources for TDS in receiving waters are agricultural and residential runoff, clayrich mountain waters, soil contamination leaching, and water pollution discharges from industrial or sewage treatment plants.

The most common chemical components are calcium, phosphates, nitrates, sodium, potassium, and chloride, which are found in nutrient runoff, general stormwater runoff, and snow runoff where road deicing salts are applied. Chemicals can be cations, anions, molecules, or aggregates of a thousand or fewer molecules, as long as soluble micro-granules are formed.

Exotic and harmful elements of TDS are pesticides resulting from surface runoff. Some naturally occurring dissolved solids occur due to the weathering and dissolution of rocks and soils. Water quality standards have been established, for example the USA has values of $500 \mathrm{mg} / \mathrm{L}$ to ensure the taste of drinking water.

The TDS value suitable for consumption is in the range of $50-100 \mathrm{ppm}$. A good classification in this sense is the following:

- Good taste: $30-60 \mathrm{ppm}$.
- Suitable taste: 60-200 ppm.
- Acceptable taste: $200-250 \mathrm{ppm}$.
- Unacceptable taste: over 250 ppm .

Total dissolved solids TDS is differentiated from total suspended solids (TSS) in that the latter cannot pass through a two micrometer sieve and yet remains suspended indefinitely in solution. The term "settleable solids" refers to material of any size that does not remain suspended or dissolved in a storage tank that is not subject to movement and excludes both TDS and TSS. Soluble solids can include larger particles or insoluble molecules.

TDS can be measured using two methods such as: Gravimetric analysis and conductometry (measurement of electrical conductivity).

Gravimetric methods are the most accurate and involve evaporating the liquid solvent and measuring the mass of residue left. This method is generally the best, although it is time consuming. If inorganic salts comprise the vast majority of TDS, gravimetric methods are appropriate.

The electrical conductivity of water is directly related to the concentration of ionized solids dissolved in the water. Ions from solids dissolved in water create the ability of that water to conduct an electrical current that can be measured using a conventional counter meter or a TDS meter. When correlated with laboratory TDS measurements, conductivity provides an approximate value for TDS concentration, usually to within ten percent.

## The relationship between TDS and the specific electrical conductance of water follows the

 following equation: $T D S=k_{e} E C$Where TDS is measured in $\mathrm{mg} / \mathrm{L}$ and EC (electrical conductance is measured in microSiemens $/ \mathrm{cm}$ at $25^{\circ} \mathrm{C}$.
$k_{e}$ is the correlation factor between the two quantities and has values between 0.55 and 0.8 .
From the point of view of TDS values, water is classified into:

- Fresh water: TDS $=500 \mathrm{ppm}$
- Brackish water: TDS $=500-30,000 \mathrm{ppm}$
- Saline water: TDS $=30,000-40,000 \mathrm{ppm}$
- Hypersaline: TDS greater than $40,000 \mathrm{ppm}$

Dissolved ionized solids such as salts and minerals increase the electrical conductivity (EC) of a solution. Because it is a measure of the volume of ionized solids, EC can be used to estimate TDS. Solid organic solids such as sugar and microscopic solid particles, or those such as colloids, do not significantly affect the conductivity of a solution and are not considered.

A TDS meter indicates the total dissolved solids (TDS) of a solution, i.e. the concentration of dissolved solid particles.

The best method for measuring the TDS of water in the laboratory is to evaporate the water and solvents and then weigh the remaining residue.
Many TDS meters display TDS values in parts per million (ppm);
1 ppm indicates 1 milligram of dissolved solids per kilogram of water.
The terms EC and TDS are often used to quantify the amount of dissolved solids in water. Pure $\mathrm{H}_{2}$ O has practically zero conductivity. The conductivity is usually about 100 times greater than the cation or anion values expressed as equivalents. TDS is calculated by converting EC by a factor of 0.5 to 1.0 times EC, depending on the levels studied.

The only accurate way to measure TDS is to evaporate the water and weigh the solid residue obtained. The TDS level can be estimated by measuring the EC (electrical conductivity) of the water with a conductometer and converting. A TDS meter is actually an EC meter calibrated in TDS, estimating TDS from EC. Some meters can be set to display either one value.

All elements dissolved in water have an electrical charge. It is therefore possible to estimate TDS values by determining the electrical conductivity of water by passing a small current through it. A similar measurement is made with an ohmmeter which measures the resulting voltage and current and calculates the conductivity.

Converting from CE to TDS can create significant errors because the conversion factor is different for different dissolved solids.

For example, for potassium chloride, to convert EC to TDS, the conversion factor is 0.50 0.57 . Sodium chloride has a factor of $0.47-0.50$ and some commonly found dissolved minerals can be as high as 0.85 .

The conversion factor also varies with temperature and only more sophisticated measurements measure solution temperature and compensate.

## Experimental assembly:

It consists of the TDS-meter, which measures the TDS values of different solutions

## PROCEDURE FOR THE PRACTICAL WORK - STEPS


1.Introduce into a Berzelius glass, the experimental solution
2. Introduce the TDS-meter into the Berzelius glass filled with solution
3. The TDS value for each solution to be studied is read on the TDSmeter display
3. Enter the obtained values in the table below

Table 1. Experimental data for TDS analysis

| No. <br> crt | The solution | TDS (ppm) | Remark |
| :--- | :--- | :--- | :--- |
| 1 |  |  |  |
| 2 |  |  |  |
| 3 |  |  |  |
| 4 |  |  |  |
| 5 |  |  |  |

## DETERMINATION OF SALINITY OF SOLUTIONS

## Theoretical introduction

Salinity is the amount of salt dissolved in an aqueous medium, called salt water. This is usually measured in $\frac{g \text { salt }}{k g \text { sea water }}$ (it is technically a dimensionless quantity).

Salinity is an important factor in determining many chemical quantities of natural waters and the biological processes within them.
Salinity is a thermodynamic state variable that, together with temperature and pressure, governs physical characteristics such as density and heat transfer capacity of water.

## Salinity in the medical field

Saline is a mixture of sodium chloride in water and has a number of uses in medicine. Applied to the affected area, it is used to clean wounds, help remove contact lenses and improve eye moisture. By injection into a vein it is used to treat dehydration, such as in cases of gastroenteritis and diabetic ketoacidosis. It is also used to dilute other drugs that must be given by injection.

Large amounts can lead to congestion, swelling, acidosis, and sodium in the blood. In those with long-term blood sodium, excessive use can cause osmotic demyelination syndrome.

Saline is part of the crystalloid family of medications. It is most commonly used as a sterile solution of 9 g of salt per liter ( $0.9 \%$ ), known as normal saline. Higher and lower concentrations may also occasionally be used.

## Saline has a pH of 5.5, which makes it acidic

The medical use of saline began around 1831.
It is on the World Health Organization's list of essential medicines, the most effective and safest medicines needed in a health system. The wholesale cost in developing countries is about $\$ 0.60-\$ 4.20$ per liter of normal saline.

Normal saline (NSS, NS, or N/S) is the usual name for a solution of $0.90 \% \mathrm{w} / \mathrm{v} \mathrm{NaCl}, 308$ $\mathrm{mOsm} / \mathrm{L}$ or 9.0 g per liter. Less commonly, this solution is referred to as physiological saline or isotonic saline (as it closely approximates isotonic, i.e., physiologically normal); although neither of these names is technically accurate (since normal saline is not exactly like blood serum), they convey
the commonly observed practical effect: good fluid balance with a minimum of hypotonicity or hypertonicity.

Normal saline is frequently used as an intravenous (IV) drip for people who cannot take fluids by mouth and have developed or are at risk of developing dehydration or hypovolemia. NS is also used for aseptic purpose. NS is typically the first fluid used when hypovolemia is severe enough to threaten the adequacy of blood circulation and has long been thought to be the safest fluid to give rapidly in large volumes.
However, it is now known that rapid infusion of NS can cause metabolic acidosis.
Saline solution is the solution of 9 grams of sodium chloride $(\mathrm{NaCl})$ dissolved in water, up to a total volume of 1000 ml (weight per unit volume (weight / volume)). The mass of 1 milliliter of normal saline is 1.0046 grams at $22^{\circ} \mathrm{C}$. The molecular weight of sodium chloride is about 58.5 grams per mole, so 58.5 grams of sodium chloride equals 1 mole. Since normal saline contains 9 grams of NaCl , the concentration is 9 grams per liter divided by 58.5 grams per mole or 0.154 moles per liter. Because NaCl dissociates into two types of ions, sodium and chloride ions - 1 molar is 2 osmolar of NaCl .

Thus, saline solution contains $154 \mathrm{mEq} / 1 \mathrm{Na}+$ and $\mathrm{Cl}-$. It has a slightly higher degree of osmolarity (ie more dissolved per liter) than blood. However, if you take into account the osmotic coefficient, a correction for non-ideal solutions, then saline is much closer to isotonic.

NaCl is about 0.93 , so $0.154 \times 1000 \times 2 \times .93=286.44$. However, the osmolarity of normal saline is close to the osmolarity of NaCl in the blood.

## 1 liter of $\mathbf{0 . 9 \%}$ saline solution contains:

- 154 mEq of sodium ion $=154 \mathrm{mmol}$
- 154 mEq of chloride ion $=154 \mathrm{mmol}$

For medical purposes, saline is often used to clean wounds and skin abrasions. Normal saline will not burn or sting when applied.

Saline is also used in intravenous therapy, the intravenous supply of supplemental water to rehydrate people or to supply the daily water and salt needs ("maintenance" needs) of a person who cannot take them orally. Because infusing a low-osmolality solution can cause problems such as hemolysis, low-concentration intravenous saline solutions typically have dextrose (glucose) added to maintain a safe osmolality while delivering less sodium chloride. The amount of normal saline infused depends largely on the person's needs (eg, diarrhea or heart failure).

Saline is often used as a nasal wash to relieve some of the symptoms of the common cold. The solution exerts a softening and loosening influence on the mucus to facilitate the washing and cleaning of the nasal passages for both babies and adults. In this case "household saline" can be used: this is done by dissolving about half a teaspoon of table salt in 240 ml of clean tap water. In very rare cases, Naegleria fowleri amoeba infection can occur if the amoeba enters the body through the nose, therefore the water used for nasal irrigation must be sterile. Note that clean tap water is not necessarily sterile liquid.

Eye drops are drops containing saline solution, used in the eye. Depending on the condition being treated, they may contain steroids, antihistamines, sympathomimetics, beta-blockers, parasympathomimetics, parasympatholytics, prostaglandins, nonsteroidal anti-inflammatory drugs (NSAIDs), antibiotics, or topical anesthetics.

Eye drops sometimes have no medicine in them and are just lubricating and tear replacement solutions.

Saline drop systems designed for syringes (eg, Wallace Cameron Ultra Saline Minipod) are distributed in modern needle exchange programs because drugs can be effectively administered either by injection or ophthalmically, which is comparable to intravenous use.

It has been shown that the elimination of latanoprost acid from plasma is rapid (half-life of 17 minutes) after ophthalmic or intravenous administration. However, ophthalmic use is done with filtered sterile drugs that are potent in fairly small doses and with pH adjusted acidity of 7.0-7.5 after drug addition to avoid eye damage. The human eye has a pH of about 7.5 , water has a pH of 7.0. Hypertonic saline solutions of $7 \% \mathrm{NaCl}$ are considered mucoactive agents and are therefore used to hydrate thick secretions (mucus) to facilitate coughing and expectoration. hypertonic saline solutions of $3 \% \mathrm{NaCl}$ are also used in cases of critical care, acutely raised intracranial pressure, or severe hyponatremia.

Inhaling hypertonic saline has been shown to help with other respiratory problems, especially bronchiolitis. Hypertonic saline is currently recommended by the Cystic Fibrosis Foundation as a primary part of a treatment regimen for cystic fibrosis.

An $11 \%$ xylitol solution with $0.65 \%$ saline stimulates nasopharyngeal lavage and has an effect on nasal pathogenic bacteria. This has been used in complementary and alternative medicine.

Other commonly used concentrations include:

- Semi-normal saline ( $0.45 \% \mathrm{NaCl}$ ), often with "D5" ( $5 \%$ dextrose), contains $77 \mathrm{mEq} / \mathrm{L}$ Na and Cl and $50 \mathrm{~g} / \mathrm{L}$ dextrose.
- $1 / 4$ normal saline $(0.22 \% \mathrm{NaCl})$ has $39 \mathrm{mEq} / \mathrm{L}$ of Na and Cl and almost always contains $5 \%$ dextrose for osmolality reasons. It can be used alone in neonatal intensive care units.

Saline can be used in perioperative fluid administration protocols to reduce excessive intravenous fluid infusions and decrease pulmonary complications.

Hypertonic saline is used to treat hyponatremia and cerebral edema. Rapid correction of hyponatremia by means of hypertonic saline or any saline infusion $>40 \mathrm{mmol} / \mathrm{L}(\mathrm{Na}+$ having a valency of $1.40 \mathrm{mmol} / \mathrm{L}=40 \mathrm{mEq} / \mathrm{L}$ ) central pontine myelinolysis $(\mathrm{CPM})$ and requires constant monitoring of the person's answer. Removal of water in combination with diuretic block does not produce the risk of CPM as high as saline administration; however, it does not correct hyponatremia as rapidly as administration of hypertonic saline. Due to hypertonicity, administration may lead to phlebitis and tissue necrosis. As such, concentrations greater than $3 \% \mathrm{NaCl}$ must normally be administered via a central venous catheter, also known as a "central line".

Such hypertonic saline is normally available in two concentrations, the first of which is more commonly administered:
$3 \% \mathrm{NaCl}$ has $513 \mathrm{mEq} / \mathrm{L}$ of Na and Cl .
$5 \% \mathrm{NaCl}$ has $856 \mathrm{mEq} / \mathrm{Na}$ and Cl .
Less commonly used NaCl solutions are $7 \%(1200 \mathrm{mEq} / \mathrm{L}$ ) and $23.4 \%$ (about $4000 \mathrm{mEq} / \mathrm{L}$ ), both of which are used (also via central line), often with additional diuretics, in the treatment of traumatic brain injuries.

Dextrose (glucose) $4 \%$ in $0.18 \%$ saline is sometimes used for maintenance.
In the medical field, the most common types of saline solutions are the soup as follows:

- Lactated Ringer's solution
- Acetate Ringer's solution
- Intravenous sugar solution
- 5\% dextrose in normal saline ( D5NS )
- $10 \%$ dextrose in normal saline solution (D10NS)
- $5 \%$ dextrose in semi-normal saline solution (D5HNS)
- $10 \%$ dextrose in semi-normal saline solution (D10HNS)

In cell biology, the following saline solutions are used:

- buffered saline (PBS) (Dulbecco recipe = D-PBS, Galfre, Kuchler, Ausubel, etc.)
- TRIS - buffered saline solution (TBS) (Goldsmith recipe, Ausubel, etc.)
- Hank's balanced salt solution (HBSS)
- Earle's balanced saline solution (EBSS)
- Standard Saline Citrate (SSC)
- HEPES - buffered saline solution (HBS) (recipes from Dittmar, Liu, Ausubel, etc.)
- Gey's balanced saline solution (GBSS)
- Lactated Ringer's solution (also known as sodium lactate solution or Hartmann's solution, is a mixture of sodium chloride, sodium lactate, potassium chloride and calcium chloride in water)


## 1 liter of lactated Ringer solution contains:

- 130-131 mEq sodium ions $=130 \underline{\mathrm{mmol} \mathrm{L}^{-1}}$
- $109-111 \mathrm{mEq}$ chlorine ions $=109 \mathrm{mmol} \mathrm{L}^{-1}$
- $28-29 \mathrm{mEq}$ of lactate ions $=28 \mathrm{mmol} \mathrm{L}^{-1}$
- $4-5 \mathrm{mEq}$ potassium ions $=4 \mathrm{mmol}^{-1}$
- $2-3 \mathrm{mEq}$ of calcium ions $=1.5 \mathrm{mmol} \mathrm{L}^{-1}$

Intravenous sugar solution is also called dextrose solution and is a mixture of dextrose and water that includes:

D5W (5\% dextrose in water), which consists of $278 \mathrm{mmol} / \mathrm{L}$ dextrose
D5NS (5\% dextrose in normal saline solution), which in addition, contains normal saline solution ( $0.9 \% \mathrm{w} / \mathrm{v} \mathrm{NaCl}$ ).

D5 1/2NS 5\% dextrose in half normal saline solution), ( $0.45 \% \mathrm{w} / \mathrm{v}$ of NaCl$)$.
D5LR (5\% dextrose in Lactated Ringer's solution)
D50-50\% dextrose in water
The percentage is here considered a mass percentage, so $5 \%$ glucose / dextrose solution contains $50 \mathrm{~g} / \mathrm{L}$ of glucose / dextrose ( $5 \mathrm{~g} / 100 \mathrm{ml}$ )

Glucose provides $4 \mathrm{kcal} / \mathrm{gram}$ of energy, so a $5 \%$ glucose solution provides $0.2 \mathrm{kcal} / \mathrm{ml}$. If prepared from dextrose monohydrate, which provides $3.4 \mathrm{kcal} / \mathrm{gram}$, a $5 \%$ solution provides $0.17 \mathrm{kcal} / \mathrm{ml}$ A salinometer is a device designed to measure the salinity or dissolved salt content of a solution.

Because salinity affects both the electrical conductivity and specific gravity of a solution, a salinometer often consists of an ec meter (conductometer) or hydrometer and some method of converting these readings into a salinity reading.

A salinometer can be calibrated either in micro-ohms, in a unit of electrical conductivity (typically 022 ), or directly calibrated in units for salinity in "grains per gallon" ( $0-0.5$ ).

A typical reading will be 2 micro-Ohms or 0.05 grains per gallon.
Fresh water generators (evaporators) use salinometers for water distillation in order to measure water quality. Evaporator water can be used for drinking water supply, salt water is not desirable for human consumption.

In some cases, higher distillation is required for use in water boilers, where salt water would be disastrous. In these cases, a salinometer is also installed on the supply system where it would alert the engineer to any salt contamination. The salinometer can automatically switch the evaporator output from the fresh water tanks to the feed water tanks depending on the water quality. Superior quality (low salinity) is required for boiler feed water, not for drinking.

A contour line of constant salinity is called an isohaline, or sometimes an isohaline.
Salinity in rivers, lakes, and the ocean is conceptually simple, but technically challenging to accurately define and measure.

Conceptually, salinity is the amount of salt dissolved in water. Salts are compounds such as sodium chloride, magnesium sulfate, potassium nitrate, and sodium bicarbonate that dissociate into ions. The concentration of dissolved chloride ions is sometimes referred to as chlorinity. Operationally, solute is defined as that which can pass through a very fine filter (a filter with a pore size of $0.45 \mu \mathrm{~m}$, but nowadays usually $0.2 \mu \mathrm{~m}$ ).

Salinity can be expressed as a mass fraction, that is, the mass of dissolved material in a unit mass of solution.

Seawater typically has a salinity value of about $35 \mathrm{~g} / \mathrm{kg}$, although lower values are typical near coasts where rivers enter the ocean.

Rivers and lakes can have a wide range of salinities, from less than $0.01 \mathrm{~g} / \mathrm{kg}$ to several $\mathrm{g} / \mathrm{kg}$, although there are many places where higher salinity values are found.

The Dead Sea has a salinity greater than $200 \mathrm{~g} / \mathrm{kg}$.
Rainwater, before it hits the ground, has a TDS of $20 \mathrm{mg} / \mathrm{l}$ or less.

Regardless of the pore size used for the filters, the resulting salinity value of a given sample of natural water will not vary by more than a few percent (\%). However, physical oceanographers working in the abyssal ocean are often concerned about the accuracy and intercomparability of measurements made by different researchers at different times to close to five significant figures. A bottled seawater product known as IAPSO Standard Seawater is used by oceanographers to standardize their measurements with sufficient accuracy to meet this requirement.

| Classification of aqueous environments according to salinity <br> levels |  |  |  |
| :--- | :--- | :--- | :--- |
| Fresh water | Brackish water | Salt water | Brine |
| $<\mathbf{0 . 0 5 \%}$ | $\mathbf{0 . 0 5 - 3 \%}$ | $\mathbf{3 - 5 \%}$ | $>\mathbf{5 \%}$ |
| $\mathbf{0 . 5 \%}$ | $\mathbf{0 . 5 - 3 0 \%}$ | $\mathbf{3 0 - 5 0 \%}$ | $>\mathbf{5 0 \%}$ |

Another name for marine waters is euhaline seas. The salinity of euhaline seas is from 30 to 35 . Seas or poor waters have salinities in the range of 0.5 to 29 and metahaline seas from 36 to 40 . These waters are considered thalassic because their ocean salinity is derived and defined as homoiohaline if salinity does not vary much over time (essentially constant). The table below modified from Por (1972), follows the "Venetia" (1959) system.

| Thalassic water series |
| :--- |
| $>300$ |
| hyperhaline |
| $60-80$ |
| metahalines |
| 40 |
| myxoeuhalines |
| 30 |
| polyhaline |
| 18 |
| mesohaline |
| 5 |
| oligohalines |
| 0.5 |



Very salty seawater, for which the salts crystallize, is also called saline solution.

Salinity is an ecological factor of considerable importance, influencing the types of organisms that live in a watery environment. Salinity also influences the types of plants that will grow either in a watery environment or on land fed by water (or groundwater). 11 A plant adapted to saline conditions is called a halophyte. A halophyte that is tolerant of residual sodium carbonate salinity is called a glass or salt or barrel plant. Organisms (mostly bacteria) that can live in very salty conditions are specifically classified as extremophiles or halophiles. An organism that can withstand a wide range of salinities is called euryhaline.

Ocean salinity is a driver of the world's ocean circulation, where density changes due to both salinity changes and ocean surface temperature changes produce changes in buoyancy, causing the water mass to sink and rise. Changes in ocean salinity are thought to contribute to global changes in carbon dioxide because more saline water is less soluble in carbon dioxide. Additionally, during glacial periods, the hydrography is such that a possible cause of reduced circulation is the production of stratified oceans. In such cases, it is more difficult to clean water through thermohaline circulation.

## Experimental assembly:

It consists of a salinometer, which measures the salinity values of different solutions

PROCEDURE FOR THE PRACTICAL WORK - STEPS

1.Introduce into a Berzelius glass, the experimental solution
2. Introduce the salinometer into the Berzelius glass filled with solution
3. The salinity value for each solution to be studied is read on the salinometer display
3. Introduce the obtained values in the table below

Table 1. Experimental data for salinity measurement

| No. <br> crt | The solution | salinity (ppt) | Remark |
| :--- | :--- | :--- | :--- |
| 1 |  |  |  |
| 2 |  |  |  |
| 3 |  |  |  |
| 4 |  |  |  |
| 5 |  |  |  |

## THE STUDY OF PH

## INTRODUCTION

SOLID ELECTROLYTES completely dissociate in aqueous solution.
Examples of strong electrolytes:
all salts,
strong acids $\left(\mathrm{H}_{2} \mathrm{SO}_{4}, \mathrm{HCl}, \mathrm{HNO}_{3}\right)$
strong bases ( $\mathrm{NaOH}, \mathrm{KOH}$ )
Weak electrolytes dissociate to a small extent. Most molecules are undissociated in aqueous solution.
Examples of weak electrolytes:
Weak acids $\left(\mathrm{H}_{2} \mathrm{CO}_{3}, \mathrm{CH}_{3} \mathrm{COOH}\right)$
Weak bases: ammonia $\left(\mathrm{NH}_{3}\right)$ and water $\left(\mathrm{H}_{2} \mathrm{O}\right)$.
Non-polar solvents can be lipids in the body, phospholipids of membranes.
They preferentially dissolve non-polar substances and not ions.
Reactions between ions in aqueous solution, at $37^{\circ} \mathrm{C}$, are fast reactions.
The acid-base systems and the systems corresponding to other ionic reactions in the body are in a state of equilibrium. This state of equilibrium may differ from one compartment to another.
The law of mass action is used to treat chemical and physical balances.
The acidity or basicity of an aqueous solution is measured by pH or pOH
The dissociation of water takes place according to the stoichiometric equation: $\mathrm{H}_{2} \mathrm{O} \leftrightarrows \mathrm{H}^{+}+\mathrm{OH}^{-}$
Or more correctly: $2 \mathrm{H}_{2} \mathrm{O} \leftrightarrows \mathrm{H}_{3} \mathrm{O}^{+}+\mathrm{OH}^{-}$
For simplification it is considered
$\mathrm{H}^{+}$the hydronium ion of the formula $\mathrm{H}_{3} \mathrm{O}^{+}$
Water dissociates very little, its dissociation constant or the ionic product of water is:
$\mathrm{H}^{+}$and $\mathrm{OH}^{-}$ions come only from the dissociation of water and will be in equal amounts:

$$
P_{\mathrm{H}_{2} \mathrm{O}}=10^{-14}(\mathrm{~mol} / \mathrm{l})^{2}
$$

In reality, water also dissolves some $\mathrm{CO}_{2}$ from the air, and some silicate from the glass, so it is not pure.

If a strong, completely dissociated acid is dissolved, practically all $\mathbf{H}^{+}$ions come from the dissociation of the acid dissolved in water. The product of the concentration of the 2 ions remains constant: $\left[\mathrm{H}^{+}\right] \cdot\left[\mathrm{OH}^{-}\right]=$constant, thus the concentration $\left[\mathrm{OH}^{-}\right]$decreases.

If a strong, completely dissociated base is dissolved, the concentration of $\mathrm{OH}^{-}$ions will be equal to the concentration of the dissolved base and the concentration of $\mathrm{H}^{+}$ions will decrease

Acidic aqueous solution $=$ is an aqueous solution where $\left[\mathbf{H}^{+}\right]>\left[\mathrm{OH}^{-}\right]$and $\mathbf{p H}<7$
Basic aqueous solution $=$ is an aqueous solution where $\left[\mathrm{H}^{+}\right]<\left[\mathrm{OH}^{-}\right]$and $\mathbf{p H}>7$
Neutral aqueous solution $=$ is an aqueous solution where $\left[\mathrm{H}^{+}\right]=\left[\mathrm{OH}^{-}\right]$and $\mathbf{p H}=7$

## DISSOCIATION OF ACIDS AND WEAK BASES

## THE NOTION OF pH

Electrolytic dissociation is a phenomenon by which some substances in solution break down into positive ions and negative ions.

By dissociating an acid, hydrogen ions and acid radical ions appear, as follows: $\mathbf{A H} \Leftrightarrow \mathbf{A}^{\boldsymbol{+}}+\mathbf{H}^{\boldsymbol{+}}$ By dissociating a base, oxidryl ions and basic radical ions are obtained, as follows: $\mathbf{B O H} \Leftrightarrow \mathbf{B}^{+}$ $+\mathrm{OH}^{-}$

Water behaves like an electrolyte and dissociates into hydrogen ions and oxydryl ions: $\mathbf{H}_{2} \mathrm{O} \Leftrightarrow$ $\mathbf{H}^{+}+\mathbf{O H}^{-}$

In reality the hydrogen ion $\mathbf{H}^{+}$obtained after dissociation hydrates and the equilibrium relationship will be written: $\mathbf{2 H O H} \Leftrightarrow \mathbf{H}_{3} \mathbf{O}^{+}+\mathbf{H O}^{-}$

## The ion $\mathrm{H}_{3} \mathrm{O}^{+}$- is called hydronium

The double-headed arrow shows that the water molecules ionize and are permanently regenerated from ions. If the temperature is constant, the speed of ionization and restoration of the molecules remains constant, achieving a stationary concentration of $\mathrm{H}^{+}$equal to the concentration of $\mathrm{OH}^{-}$. If a strong acid is dissolved in water, the concentration of hydronium ions $\mathrm{H}^{+}$increases and that of hydroxyl ions $\mathrm{OH}^{-}$decreases, because the ionic product of water remains constant at a given temperature.If a strong base is dissolved in water, the concentration of $\mathrm{OH}^{-}$increases and that of $\mathrm{H}^{+}$ decreases.

At constant temperature the concentrations of $\mathrm{H}^{+}$and $\mathrm{OH}^{-}$will oscillate as follows:

| $\mathbf{H}^{+}$ | $10^{0}$ | $10^{-1}$ | $10^{-2}$ | $\cdot$ | $10^{-7}$ | . | $10^{-12}$ | $10^{-13}$ | $10^{-14}$ |
| :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- |
| $\mathbf{O H}^{-}$ | $10^{-14}$ | $10^{-13}$ | $10^{-12}$ | $\cdot$ | $10^{-7}$ | . | $10^{2}$ | $10^{1}$ | $10^{0}$ |
| Character | ACID |  |  |  |  |  |  |  |  |

To simplify the calculations, it was agreed to express this concentration
by $\mathbf{p H}$ - which represents the logarithm with changed sign of the hydrogen ion concentration:
$\mathbf{p H}=-\boldsymbol{\operatorname { l o g }} \mathrm{H}^{+}$
or in general, according to SORENSEN: $\mathrm{pX}=-\log \mathrm{X}=\log (1 / \mathrm{X})$
also, it can be written as follows: $\mathrm{pH}+\mathrm{pOH}=\mathrm{pK}$, where pK is the ionic product of water

| pH | 0 | 1 | 2 | $\cdot$ | 7 | . | 12 | 13 | 14 |
| :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- |
| pOH | 14 | 13 | 12 | . | 7 | . | 2 | 1 | 0 |
| Character | ACIDITY | NEUTRALITY | ALKALINE |  |  |  |  |  |  |

The acid or basic character of a solution can be described quantitatively in terms of the concentration of hydrogen ions. In some situations, the concentration has a very low value, thus it is more convenient to use the notion of pH .

The $\mathbf{p H}$ is by definition the decimal logarithm with changed sign of the molar concentration of hydronium ions:
$p H=-\log \left[H^{+}\right]$

## The pOH of a solution is:

$\mathrm{pOH}=-\log \left[O H^{-}\right]$
The relationship between them is: $\mathrm{pH}+\mathrm{pOH}=14$
pH and pOH are complementary

It is enough to indicate the pH to know if a solution is acidic or alkaline:

- if $\mathbf{p H}<\mathbf{7}$ - the solution is acidic
- if $\mathbf{p H}=7$ - the solution is neutral
- if $\mathbf{p H}>7$ - the solution is alkaline (or basic)
pH SCALE

| $\mathbf{H}+$ | $\mathbf{1 0}^{\mathbf{0}}$ | - | $\mathbf{1 0}^{-\mathbf{7}}$ | - | $\mathbf{1 0}^{-\mathbf{1 4}}$ |
| :--- | :--- | :--- | :--- | :--- | :--- |
| pH | 0 | - | 7 | - | 14 |
| pOH | 14 | - | 7 | - | 0 |
| Character | ACID | NEUTRAL | BASIC |  |  |

The neutral solution is the solution in which the concentration of $\mathrm{H}^{+}$and $\mathrm{OH}^{-}$ions are equal and consequently the value of $\mathbf{p H}$ is $\mathbf{7}$.
When a strong acid is dissolved in water, it completely dissociates into ions, the $\mathrm{H}^{+}$ions practically only come from the ionization of the acid. The ionic product of water remains constant, the increase in the concentration of $\mathrm{H}^{+}$ions remains constant, it will cause the decrease in the concentration of $\mathrm{OH}^{-}$ions,
the acid range is the one in which $\mathbf{p H}<7$.
If a strong base is dissolved in water, $\mathbf{O H}^{-}$ions will come exclusively from this ionization, the concentration of $\mathrm{OH}^{-}$ions will prevail.
The basic domain is the one in which the $\mathbf{O H}^{-}$concentration is higher, thus: $\mathbf{p H}>7$.

ELECTROLYTES AND NON-ELECTROLYTES

| ACID |  | BASES |  |
| :---: | :---: | :---: | :---: |
| Strong acids | Weak acids | Strong bases | Weak bases |
| SOLID INORGANIC ACIDS <br> nitric acid, $\mathrm{HNO}_{3}$ hydrobromic acid, HBr hydrochloric acid, HCl hydrochloric acid, $\mathrm{HClO}_{3}$ hydroiodic acid, HI perchloric acid, $\mathrm{HClO}_{4}$ sulfuric acid, $\mathrm{H}_{2} \mathrm{SO}_{4}$ | WEAK <br> AND <br> MEDIUM <br> STRENGTH <br> INORGANIC <br> ACIDS <br> boric acid, $\mathrm{H}_{3} \mathrm{BO}_{3}$ <br> carbonic acid, $\mathrm{H}_{2} \mathrm{CO}_{3}$ (simply called <br> "acid" when included in mineral <br> waters or soft drinks such as Fanta) <br> hydrofluoric acid, HF <br> phosphoric acid $\mathrm{H}_{3} \mathrm{PO}_{4}$ <br> WEAK ORGANIC ACIDS <br> acetic acid, $\mathrm{CH}_{3} \mathrm{COOH}$ <br> butyric acid, $\mathrm{CH}_{3} \mathrm{CH}_{2} \mathrm{CH}_{2} \mathrm{COOH}$ <br> benzoic acid, $\mathrm{C}_{6} \mathrm{H}_{5} \mathrm{COOH}$ <br> citric acid, $\left(\mathrm{C}_{6} \mathrm{O}_{7} \mathrm{H}_{8}\right)$ <br> formic acid, HCOOH <br> lactic acid, $\mathrm{CH}_{3} \mathrm{CHOHCOOH}$ <br> malic <br> acid, <br> $\mathrm{HOOCCH}_{2} \mathrm{CHOHCOOH}^{2}$ <br> oxalic acid, HOOCCOOH <br> pyruvic acid, CHCOCOOH <br> propionic acid, $\mathrm{CH}_{3} \mathrm{CH}_{2} \mathrm{COOH}$ <br> valerianic acid, $\mathrm{CH}_{3}\left(\mathrm{CH}_{2}\right)_{3} \mathrm{COOH}$ | NaOH <br> KOH <br> LiOH <br> $\mathrm{Ca}(\mathrm{OH})_{2}$ <br> $\mathrm{Ba}(\mathrm{OH})_{2}$ | Ammonia <br> $\mathrm{NH}_{3}$ <br> $\mathrm{H}_{2} \mathrm{O}$ <br> $\mathrm{Mg}(\mathrm{OH})_{2}$ <br> $\mathrm{Zn}(\mathrm{OH})_{2}$ <br> $\mathrm{Al}(\mathrm{OH})_{3}$ |

## ELECTROLYTES AND NON-ELECTROLYTES

## calculation formulas for $\mathbf{p H}$

- For STRONG ACIDS : $p H=-\lg \left[H^{+}\right]=-\lg [H A]=-\lg c$
- For WEAK ACIDS : $p H=p K_{H A}+\lg \frac{[D]}{[P]}$
(Henderson Hasselbach formula )
D-form deprotonated, P- form protonated
- For STRONG BASES: $p H=14-p O H=14+\lg \left[O H^{-}\right]$


## The pH of SOLUTIONS OF WEAK ACIDS AND BASES

- When only the initial concentration of the weak acid is known, taking into account that the $\mathrm{H}^{+}$and $\mathrm{A}^{- \text {ions }}$ resulting from the dissociation are in the solution in equal concentrations, the equilibrium constant becomes: $K_{H A}=\frac{\left[H^{+}\right] .{ }^{2}}{[H A]}$
- After the logarithm, we get: $\left.{ }^{-\lg K} H A=-\lg \left[H^{+}\right]\right]^{2}+\lg [H A]$
- Similar to the definition of $\mathbf{p H}: \mathbf{p H}=-\lg K$
- The formula for calculating $\mathbf{p H}$ is:

$$
p H=\frac{p K_{H A}-\lg [H A]}{2}
$$

When the concentration of the protonated form is known
$[\mathbf{P}]=[\mathbf{H A}]-$ the molar concentration of the weak acid
as well as the concentration of the deprotonated form $[\mathbf{D}]=\left[\mathbf{A}^{-}\right]$- the molar concentration of the conjugate base

- The HENDERSON HASSELBACH FORMULA $p H=p K_{H A}+\lg \frac{[D]}{[P]}$ is used :

EXPERIMENTAL ASSEMBLY: pH - meter


1- keyboard (Keypad)
2- display
3- the socket system (Jack field)


Affichage

|  | affichage de l'état |
| :---: | :---: |
|  | symbole de senseur / l'électrode affichage de la valeur mesurée |
|  | fonction et affichage de la température |

domaine de prise


| $\mathbf{1}$ | D.O. sonde ou la conductivité du cellule |
| :--- | :--- |
| $\mathbf{2}$ | Electrode de pH |
| $\mathbf{3}$ | Sonde de température de pH |
| 4 | Branchez - dans l'unité d'alimentation |
| $\mathbf{5}$ | RS232 interface série |

## PROCEDURE FOR THE PRACTICAL WORK - STEPS



## DETERMINATION IN THE ACID MEDIUM:

1. Firstly, the calibration is performed in the acid range by using a buffered solution with $\mathbf{p H}=$ 4
2. Insert the electrode into the calibration solution and read the pH . If the device does not indicate correctly, calibrate the device.
3. Unbuffered solution (tap water) is introduced into a Berzelius glass ( $\mathrm{V}=20 \mathrm{~mL}$ )
4. Add 0.5 mL of 0.1 M HCl 5 times, by using a pipette
5. The solution is stirred until the pH -meter does not indicate any variation of pH value
6. Read the pH values and introduce them in the table
7. Repeat steps 3-6 with an unbuffered solution of $\mathrm{pH}=4$

## DETERMINATION IN THE BASIC MEDIUM:

1. Firstly, the calibration is made in the basic range using the buffered solution with $\mathbf{p H}=\mathbf{7}$
2. Insert the electrode into the calibration solution and read the pH . If the device does not indicate correctly, calibrate the device.
3. Unbuffered solution (tap water) is introduced into a Berzelius glass ( $\mathrm{V}=20 \mathrm{~mL}$ )
4. Add 5 times 0.5 mL of KOH 0.1 M , by using a pipette
5. The solution is stirred until the pH -meter does not indicate any variation
6. Read the pH values and introduce them in the table
7. Repeat steps 3-6 with an unbuffered solution of $\mathrm{pH}=7$

## Preparation of buffered solution with $\mathbf{p H}=4.0$

7.71 mL of $\mathrm{Na}_{2} \mathrm{HPO}_{4}$ solution of $0.2 \mathrm{M}+12.29 \mathrm{~mL}$ of citric acid solution $\left(\mathrm{C}_{6} \mathrm{H}_{8} \mathrm{O}_{7}\right)$ with 0.1 M results in 20 mL of buffered solution with $\mathrm{pH}=4$

## Preparation of buffered solution with $\mathbf{p H}=5.0$

$10.30 \mathrm{~mL} \mathrm{Na}_{2} \mathrm{HPO}_{4}$ solution of concentration $0.2 \mathrm{M}+9.70 \mathrm{~mL}$ citric acid solution $\left(\mathrm{C}_{6} \mathrm{H}_{8} \mathrm{O}_{7}\right)$ with 0.1 M we obtain 20 mL buffered solution with $\mathrm{pH}=7$.

We have a buffered solution with $\mathrm{pH}=4$ in which we introduce 5 x 0.2 mL of 0.1 M HCl and in the buffered solution with $\mathrm{pH}=7$ we introduce 5 x 0.2 mL of KOH 0.1 M .

Table 1: pH values for the basic and acid medium $\mathrm{V}=20 \mathrm{~mL}$

| No. | V (ml) | $\mathbf{p H}$ <br> acid domain |  | $\mathbf{p H}$ <br> basic domain |  |
| :---: | :---: | :---: | :---: | :---: | :---: |
|  |  | 0.1 M HCl | 0.1 M HCl | KOH 0.1M | KOH 0.1M |
|  |  | Solution unbuffered | Solution buffered | Solution unbuffered | Solution buffered |
| 1 | 0 |  |  |  |  |
| 2 | 0.5 |  |  |  |  |
| 3 | 1.0 |  |  |  |  |
| 4 | 1.5 |  |  |  |  |
| 5 | 2.0 |  |  |  |  |
| 6 | 2.5 |  |  |  |  |

## CALCULATIONS AND GRAPHICS

1. For all 4 cases, the following calculations are made:

The pH variation for the case of a solution without HCl or KOH
The pH variation for the case of a solution with HCl or KOH
$\Delta \mathrm{pH}=\mathrm{pH}($ after 25 ml$)-\mathrm{pH}(0)$
$\Delta \mathrm{pH}=$ final $\mathrm{pH}-$ initial pH
2. It is calculated $i=\frac{\Delta V}{|\Delta p H|}$ - the buffering capacity that characterizes the effectiveness of the

## buffer system to resist $\mathbf{p H}$ variations

3. It is also calculated for all 4 situations:
4. acid domain - unbuffered solution
5. acid domain - buffered solution
6. basic domain - unbuffered solution
7. basic domain - buffered solution
8. Plot the graphs $\mathrm{pH}=\mathrm{pH}(\mathrm{V})$ : for the 4 situations:
9. acid domain - unbuffered solution
10. acid domain - buffered solution
11. basic domain - unbuffered solution
12. basic domain - buffered solution

Table 2: quantities required to obtain buffer solutions with different pH

| pH | $\mathrm{Na}_{2} \mathrm{HPO}_{4}(\mathrm{c}=0.2 \mathrm{M})$ <br> $(\mathrm{mL})$ | $\mathrm{C}_{6} \mathrm{H}_{8} \mathrm{O}_{7}(\mathrm{c}=0.1 \mathrm{M})$ <br> $($ citric acid $)(\mathrm{mL})$ |
| :--- | :--- | :--- |
| 3.0 | 20.55 | 79.45 |
| 4.0 | 38.55 | 61.45 |
| 5.0 | 51.5 | 48.5 |
| 6.0 | 63.15 | 36.85 |
| 7.0 | 82.35 | 17.65 |
| 8.0 | 97.25 | 2.75 |

Table 3: quantities required to obtain buffer solutions with different pH

| pH | $\mathrm{CH}_{3} \mathrm{COOH}$ <br> $(\mathrm{c}=0.2 \mathrm{M})(\mathrm{mL})$ | $\mathrm{CH}_{3} \mathrm{COONa}$ <br> $(\mathrm{c}=0.2 \mathrm{M})(\mathrm{mL})$ |
| :--- | :--- | :--- |
| 3.47 | 95 | 5 |
| 3.80 | 90 | 10 |
| 4.15 | 80 | 20 |
| 4.38 | 70 | 30 |
| 4.57 | 60 | 40 |
| 4.75 | 50 | 50 |
| 4.93 | 40 | 60 |
| 5.12 | 30 | 70 |
| 5.70 | 10 | 90 |
| 6.03 | 5 | 95 |

## PHYSICAL TITRATION OF SOME SALTS CONDUCTOMETRICAL TITRATION

## THEORETICAL INTRODUCTION

According to the following sources: DEX '98 (1998) ; DOOM 2 (2005) ; NODEX (2002)
TITRATION has the following meaning: To determine by volumetric analysis the concentration of a solution. - From brother titer. To subject to special analysis (to determine concentration). $/<$ fr. titer To determine the proportion of components in a substance. [< bro. title ].

## TITRATION represents the quantitative analytical determination of a sample.

## TITRATION - THE GENERAL CASE

## Classification of electroanalytical methods

Electrochemical methods can be classified according to several criteria:

- according to the way in which the studied system is disturbed, they can be stationary and nonstationary;
- according to the monitored variable: current (chronoamperometry, voltammetry), potential (chronopotentiometry), load (chronocoulometry), impedance.
- according to the disturbed variable there can be: controlled potential (potentiostatic), controlled current (galvanostatic) and coulostatic methods;
- A separate category is the hydrodynamic methods (voltammetry on disk electrode or rotating ring), in which the working electrode and the electrolyte are in relative motion to each other. In these methods, forced convection works.


## Conductometric titration - principle

$\square$ the presence of ions in a solution allows the passage of electric current
$\square$ conductivity
$\square$ varies proportionally with the concentration of the analyte
$\square$ is the inverse of electrical resistance
$G=\frac{1}{R} ; \quad R=\frac{\rho \cdot l}{S} ; \quad G_{s}=\frac{1}{\rho} ; \quad G=\frac{1}{R}=G_{s} \cdot \frac{S}{l}$
$\mathrm{G}=$ conductivity
$\rho=$ resistivity
$1=$ column length
$\mathrm{S}=$ cross-sectional area of the column
G s $=$ specific conductivity

## CONDUCTOMETRICAL TITRATION - TITRATION CELL - THEORY

the glass cell, closed or not, into which the platinum or silver electrodes and the solution to be analyzed are inserted;a constant voltage is applied between the two electrodes, determining the passing current by solution as a result of a chemical reaction;
$\square$ the potential difference between the electrodes depends on the processes that take place at the electrodes and current transport phenomena.

## CONDUCTOMETRICAL TITRATION - PRINCIPLE SCHEME


$\square$ measuring the resistivity of the solution with a Wheatstone bridge
$\square$ cell constant $1 / \mathrm{s}$ ratio
conductivity specifies the ratio $\chi=\mathrm{C} / \mathrm{R}$
conductivity measurement unit - Siemens (S)

## CONDUCTOMETRICAL TITRATION - APPLICATIONS

- MONITORING THE PURITY OF DISTILLED WATER IT WAS THE BASIS OF THE DETERMINATION OF THE IONIC PRODUCT OF WATER
- DETERMINATION OF EQUIVALENCE POINTS IN TITRATIONS
- TITRATION OF STRONG ACIDS WITH STRONG BASES

A- Conductometric titration curve of 0.1 M HCl with 0.1 M NaOH
A - they have a decrease in the conductivity of the solution up to the equivalence, as a result of the formation of water molecules and the replacement of the hydrogen ion with sodium with a lower mobility. After the equivalence point, there is an increase in the conductivity of the solution as a result of the excess sodium hydroxide added and the high mobility of the hydroxyl ion.
B- Titration curve of 0.001 M pnitrophenol with 0.1 M NaOH
B - a different shape of the titration curve, which is explained by the slow increase in conductivity of the solution by the formation of ionized sodium p-nitrophenoxide, and after equivalence, a more pronounced increase in the conductivity of the solution is observed, as a result of the excess sodium hydroxide added, which is a strong electrolyte.

## - TITRATION OF WEAK ACIDS WITH WEAK BASES

In the case of the conductometric titration of a 0.01 M ammonia solution with a 0.5 M acetic acid solution in methanol (both electrolytes being weak), a curve with a totally different appearance is obtained from the other two .

Up to equivalence, the conductivity of the solution increases a lot due to the formation of ionized ammonium acetate (hard electrolyte), and after equivalence the conductivity created very slightly due to the small excess of added acetic acid solution, which is a weak electrolyte. The method based on following the variation of the conductivity of a solution as a result of a chemical reaction;

It uses a specific cell provided with two Pt electrodes, through them a constant voltage is applied determining the current passing through the solution as a result of a chemical reaction;

Through the graphic representation of the conductivity of the solution at different moments of the titration depending on the volume of titrant added, the titration curve is obtained consisting of two right segments whose intersection indicates the equivalence point.

## - POTENTIOMETRIC TITRATION

The variation of the potential of the solution to be analyzed determined by a certain chemical reaction is monitored;
$\square$ From this variation determined with the help of an indicator electrode (in relation to a reference electrode), the analyte concentration can be determined, respectively the titrant volume (from which the analyte concentration can be calculated);

## POTENTIOMETRIC TITRATION - INDICATOR ELECTRODES

## TYPE I ELECTRODES:

Metal in contact with its cation or a non-metal in contact with its anion
$\square$ Examples:

- $\quad \square \mathrm{Cu} 2+/ \mathrm{Cu}(\mathrm{s})$
- $\quad \mathrm{Zn} 2+/ \mathrm{Zn}(\mathrm{s})$
- $\square$ SHE
- $\square \mathrm{Ag}+/ \mathrm{Ag}$
- $\square \mathrm{Cl}-/ \mathrm{Cl} 2(\mathrm{~g}) / \mathrm{Pt}$
- $\quad \square$ The response of the electrode is Nernstian
- $\mathrm{E}=\mathbf{E} 0+(\mathrm{RT} / \mathrm{nF}) \ln \mathbf{a}(\mathrm{M} 2+)$


## Type II electrodes:

- Metal in contact with a sparingly soluble metal salt
- Examples:
- $\mathrm{Ag} / \mathrm{AgCl}(\mathrm{s})$
- $\mathrm{Hg} / \mathrm{Hg} 2 \mathrm{Cl} 2(\mathrm{~s}) / \mathrm{Cl}$ - (saturated calomel electrode; SCE)
- The response of the electrode is Nernstian
- $\mathrm{E}=\mathrm{E} 0-(\mathrm{RT} / \mathrm{F}) \ln \mathrm{a}(\mathrm{Cl}-)$


## Type III electrodes:

Electrodes that serve as a source of electrons (redox, inert electrodes)
Examples:

- metals: $\mathrm{Pt}, \mathrm{Au}$, graphite, Hg
- semiconductors: $\mathrm{Si}, \mathrm{GaAs}, \mathrm{In}-\mathrm{SnO} 2 /$ glass
- Answer:
- for Pt in contact with $\mathrm{Fe} 2+, \mathrm{Fe} 3+$ in solution:
- $\mathrm{E}=\mathrm{E} 0-0.059 \log ([\mathrm{Fe} 2+] /[\mathrm{Fe} 3+])$


## Potentiometric titration - applications

## Example of potentiometric titration of a chloride with AgNO3:

- the determination of the equivalence point can be done by the graphic representation of of the potential of the solution during the titration depending on the volume of added nitrate (a); received
the derivative of the potential
(b) or the second derivative of the potential
(c) depending on the volume of titrant added.

The equivalence volume can be determined from

- graphic method
- calculations based on the Hosteller - Roberts relationship. Constant volumes of titrant of 0.1 or 0.2 ml are added around the equivalence point, therefore $\Delta \mathrm{V}$ and $\Delta \mathrm{V}^{2}$ are constant and E or $\Delta \mathrm{E}^{2}$ can be graphically represented $\Delta$ as a function of V .
- $V_{e}=V_{1}+0.1 .\left(\frac{\frac{\Delta^{2} E_{1}}{\Delta V^{2}}}{\frac{\Delta^{2} E_{1}}{\Delta V^{2}}-\frac{\Delta^{2} E_{21}}{\Delta V^{2}}}\right)=V_{1}+0.1\left(\frac{\Delta^{2} E_{1}}{\Delta^{2} E_{1}-\Delta^{2} E_{2}}\right)$

Old - vol of equivalence
V1 vol of titrant before equivalence
0.1 - volume of titrant around the equivalence point
$\Delta \mathrm{E}_{2}{ }^{2} / \Delta \mathrm{V}^{2}$ - second derivative after equivalence
$\Delta \mathrm{E}^{2}{ }_{1} / \Delta \mathrm{V}^{2}$ - second derivative before equivalence

## EXPERIMENTAL ASSEMBLY



- 1 - conductometer
- 2 - burette
- 3 - thermometer
- 4 - indicator electrode - conductometric sensor
- 5 - magnetic stirrer


## PROCEDURE FOR THE PRACTICAL WORK - STEPS

A sample is introduced into a vessel whose quantity is unknow. A solvent is introduced over the sample - in our case it can also be water, after which a titration reagent is introduced under control conditions and in small portions.

During this time, properties of the contents of the vessel are monitored, such as:

- Chemical properties: chemical titration
- Physical properties: physical titration: conductivity, concentration, or optical properties such as the transparency of the solution used or not

If the titration is continued above a certain limit value, a drastic change takes place.
We want to determine the amount of $\mathrm{AgNO}_{3}$ in a sample.
The reaction takes place: $\mathrm{AgNO}_{3} \longrightarrow \mathrm{Ag}^{+}+\mathrm{NO}_{3}{ }^{-}$
The titration solution used is KCl 0.1 M which is added in small and controlled portions
The reaction takes place: $\mathrm{Ag}^{+}+\mathrm{Cl}^{-} \quad \underset{\sim}{~} \mathrm{Cl}$ (white precipitate - very difficult to dissolve solution)
The titrated sample is approximately 220 mg of $\mathrm{AgNO}_{3}$ in $\mathrm{H}_{2} \mathrm{O}$
The titration reagent is: 0.1 M KCl solution.
The titration takes place in a titration vessel (with 3 inlets) into which:

- the conductivity electrode
- the tip of the burette through which the titration reagent is added
- thermometer for temperature control.


## WORKING STEPS

1. The sample of AgNO 3 - a crystalline solid - is introduced into the titration vessel
2. An appropriate volume of water is added to bring the level of the liquid in the titration vessel to the appropriate height so that the electrode can perform the measurement.
3. The contents of the titration vessel are stirred using a magnetic stirrer.
4. The buret with the nominal volume of 25 ml is filled with the titration solution $(\mathrm{KCl} 0.1 \mathrm{M})$, up to division 0 (top).
5. While the contents of the titration vessel are being stirred, the titration solution is added in portions of 0.5 ml .
6. After each added amount, read the conductance of the solution in the titration vessel.
7. The data are entered in the table below.
8. If in the titration vessel at the end of the experiment the color of the solution obtained in the vessel is milky white, upon the introduction of ammonia the solution becomes transparent.

Table 1. The conductance of the solution in the titration vessel

| No. | $\mathrm{V}(\mathrm{ml})$ | $\mathrm{G}(\mu \mathrm{S} / \mathrm{cm})$ |
| :--- | :--- | :--- |
| 1 | 0 |  |
| 2 | 0.5 |  |
| 3 | 1 |  |
| 4 | 1.5 |  |
| 5 | 2 |  |
| $\cdots$ | $\cdots \ldots . . . .$. |  |
|  |  |  |
|  |  |  |
|  |  |  |
|  |  |  |
|  |  |  |
|  |  |  |
|  |  |  |
|  |  |  |
| 36 | 17 |  |

## EXPERIMENTAL DATA PROCESSING

At the beginning of the titration, the electrical conductivity shows a
decrease with the addition of the titration reagent because in this phase the conductivity is controlled by the very weak solubility and dissociation of AgCl . From a certain added volume, an increase in conductivity is found. Explication it consists in the fact that when the entire amount of $\mathrm{Ag}+$ is precipitated, when the titration solution is added, the total concentration of ions progressively increases ( $\mathrm{K}+$ and $\mathrm{Cl}-$ : both types of ions contribute to electrical conductivity)

The conductivity measured according to the volume of added solution is represented graphically.
The volume of the solution at which the sudden variation of the slope of the titration curve occurs is called the equivalence volume from which the amount of AgNO 3 in the sample can be calculated.

If we know the molar mass of $\mathrm{AgNO}_{3}: \mathrm{M}_{\mathrm{AgNO}}=169.86 \mathrm{~g} / \mathrm{mol}$.
$(\mathrm{MAg}=107.87 \mathrm{~g} / \mathrm{mol}, \mathrm{MN}=14 \mathrm{~g} / \mathrm{mol}, \mathrm{MO}=16 \mathrm{~g} / \mathrm{mol})$
KCl solution 0.1 M

1000 ml soil .....( contains) ............ 0.1 mol KCl
V (ml)
n mol KCl n mol AgNO3
$n=\left(V \cdot \frac{0.1}{1000}\right) \mathrm{mol} \mathrm{KCl}$
$n=\left(V \cdot \frac{0.1}{1000}\right) \mathrm{mol} \mathrm{AgNO}_{3}$
$\mathrm{M}_{\mathrm{AgNO} 3}=169.86 \mathrm{~g} / \mathrm{mol}$
$\mathrm{M}_{\mathrm{AgNO}}=169.86 \mathrm{~g} / \mathrm{mol}$
1 mole of AgNO 3


$$
\begin{aligned}
& m=\left(M_{A g N O 3} \cdot V \cdot \frac{0.1}{1000}\right) g \mathrm{AgNO}_{3} \\
& m=\left(M_{A g N O 3} \cdot V \cdot \frac{0.1}{1000}\right) g A g^{+}
\end{aligned}
$$

The conductivity measured as a function of the added solution volume is graphically represented: $\mathrm{G}=\mathrm{G}(\mathrm{V})$

It is obtained: $\mathrm{G}=\mathrm{A} . \mathrm{V}+\mathrm{B}$
Where A represents the slope of the line.
But we are dealing with two rights:
First right up to $\mathrm{Ve}: \mathrm{G}=\mathrm{A}_{1} . \mathrm{V}+\mathrm{B}{ }_{1}$
Second right from Ve up: $G=A_{2} . V+B{ }_{2}$
$\mathrm{A}_{1}$ and $\mathrm{A}_{2}$ are the slopes of the graphed lines.
To find the equivalent volume, the equations of the two lines are equal:
$\mathrm{A}_{1} \cdot \mathrm{Ve}+\mathrm{B}_{1}=\mathrm{A}_{2} \cdot \mathrm{Ve}+\mathrm{B}_{2}$
$\mathrm{Ve}\left(\mathrm{A}_{1}-\mathrm{A}_{2}\right)=\mathrm{B}_{2}-\mathrm{B}_{1}$
$V_{e}=\frac{B_{2}-B_{1}}{A_{1}-A_{2}}$

## RESISTANCES STUDY

## THEORETICAL INTRODUCTION

OHM'S LAW FOR A CIRCUIT PORTION

Considering an electric circuit composed of several consumers and a power generator, we can easily appreciate that the electric potential between points A and B decreases when the electric current passes. Carrying out the ratio of the electrical voltages $U$ and the currents I passing through each consumer, it is found that they are constant.
$\frac{U_{1}}{I_{1}}=\frac{U_{2}}{I_{21}}=\frac{U_{3}}{I_{3}}=\ldots=\frac{U_{n}}{I_{n}}=$ const .
ou
$\frac{U}{I}=R$


The ohmic conductor is an electronic component also called resistance that verifies Ohm's law:
$U=I . R$
Where I is the intensity of the current, in amperes, that crosses the resistance R , and U is the voltage at the terminals.


The characteristic curve of the intensity of the electric current I depending on the voltage at the terminals U is a straight line that passes through the origin

The characteristic of an ideal resistance is the curve: $I=f(U)$ $=U / R$

Sometimes the term pure resistance or ideal resistance is used.

For example, the resistance of a metal conductor that depends on the temperature $\theta$ is given by the relation:
$R=R_{0}\left(1+a \theta+b \theta^{2}\right)$ with $R_{0 a}$ hypothetical ohmic conductor that perfectly models a thermostated conductor at a temperature of 0 K . $\theta$ is expressed in K

The voltage drop U on a portion of the circuit is proportional to the intensity of the current I that crosses the portion of the circuit $\mathrm{U}=\mathrm{RI}$
The proportionality constant between voltage and electric current is called electrical resistance and is denoted by R. The unit of measure for electrical resistance is derived as follows.
$\frac{U}{I}=R$
$[R]_{S I}=\frac{[U]_{S I}}{[I]_{S I}}=\frac{1 V}{1 A}=1 \Omega$
Electrical resistance is characteristic of an electrical consumer and depends on its constructive elements: $R=\frac{\rho \cdot l}{S}$

Where pis a characteristic of the material from which the consumers are made.
Electrical resistivity depends on the temperature of the conductor $\rho=\rho_{0}(1+\alpha . t)$
Where $\rho o$ is the resistivity at 0 oC and $\alpha$ is the thermal coefficient of resistivity
The electrical resistance depends on the temperature $R=R{ }_{0}(1+\alpha . t)$

Table with the electrical characteristics of some materials
(http://msabau.xhost.ro/?Fizic\�:Electrokinetics:Ohm's_Law)

| the substance | $\rho_{0}(\Omega . \mathrm{m})$. <br> $10^{-8}$ | $\rho(\Omega . \mathrm{m})$ <br> $\left(20^{0} \mathrm{C}\right) .10^{-8}$ | $\alpha\left(\right.$ degree $\left.^{-1}\right)$. <br> $10^{-4}$ |
| :--- | :--- | :--- | :--- |
| nickel | 30 | 42 | 1 |
| gold | 1.92 | 2.24 | 83 |
| copper | 1.48 | 1.68 | 68 |
| iron | 8.59 | 9.71 | 65 |
| silver | 1.42 | 1.59 | 61 |
| tungsten | 5.02 | 5.47 | 45 |
| platinum | 9.83 | 10.6 | 39 |
| aluminum | 2.44 | 2.65 | 43 |
| Mercury | 94.1 | 95.8 | 9 |

## OHM'S LAW FOR AN INTEGRATED CIRCUIT

Let there be a simple electric circuit consisting of a generator that has electromotive voltage $E$ and internal resistance r , which generates an electric current of intensity I in the circuit formed by the generator and an electric resistor of resistance R , it can be written

- Ohm's law for the whole circuit : $I=\frac{E}{R+r}$
- Ohm's law for a portion of a circuit is written:

$$
\mathbf{U}=\mathbf{R I}, \mathbf{u}=\mathbf{R} \mathbf{I}
$$

So we can write the intensity of the electric current that crosses a closed electric circuit is directly proportional to the electromotive voltage E of the source and inversely proportional to the total electric resistance of the circuit.

The voltage at the terminals of the closed circuit is: $\mathrm{U}=\mathrm{E}-\mathrm{RI}$
The voltage at the terminals of the open circuit is: $\mathrm{U}=\mathrm{E}$, because $\mathrm{I}=0$
For a zero external resistance, the intensity of the short-circuit current can be written as Icc $=\mathrm{E} / \mathrm{r}$. The short-circuit current is the maximum current that an electric generator can supply

## KIRCHHOFF'S LAWS

Kirchhoff's laws express the conservation of energy and charge in an electrical circuit. They are named after the German physicist who introduced them in 1845: Gustav Kirchhoff. In a complex circuit, it is possible to calculate the potential differences across each resistor and the forward current in each circuit branch by applying Kirchhoff's two laws: the node law and the loop law.

## THE LAW OF KNOTS

The sum of the intensities of the currents entering through a node is equal to the sum of the intensities of the currents leaving the same node.

The figure shows the direction of orientation of each wire (chosen arbitrarily), which enters or leaves node A . The intensity of a current is an algebraic quantity (positive or negative) defined according to the orientation of the wire. For example, if the intensity of an input wire is 3 A , it means that in this wire at the output there will be an intensity of - 3 A .

According to the law of knots, it can be written: $\mathrm{i} 1+\mathrm{i} 2=\mathrm{i} 3+\mathrm{i} 4$.


The diagram illustrating the law of knots

This law derives directly from the law of conservation of electric charge, taking into account the fact that, in stationary mode, these charges cannot accumulate at any point in the circuit. Tasks arriving at a node offset those leaving. This law allows the solution of the "electrical equations" with the help of the knot method.

## THE LAW OF MESHES

In any EYE of a network, near the quasi-steady regimes and provided that the variations of the magnetic flux through the meshes are negligible, the algebraic sum of the potential differences along the mesh is constantly zero.

This law results from the additivity of potential differences between two points.
The potential difference between a and b is $\mathrm{Uab}=\mathrm{Va}-\mathrm{Vb}$.
Va and Vb being the respective potentials from points a and b .
By adding all these differences on a closed eye, we get a null result.


This law allows solving the electrical equations using the mesh method

```
\(\mathrm{E}_{1}+\mathrm{E}_{2}-\mathrm{E}_{3}-\mathrm{E}_{4}=\mathrm{R}_{1} \mathrm{I}_{1}-\mathrm{R}_{2} \mathrm{I}_{2}-\mathrm{R}_{3} \mathrm{I}_{3}-\mathrm{R}_{4} \mathrm{I}_{3}+\mathrm{R}_{5} \mathrm{I}_{4}\)
```

The knot law is valid only if the flux of the electric field surrounding each knot remains zero or constant. Therefore, it is not valid in electrostatics.

In a semiconductor, two types of mobile charges coexist, electrons and holes. The total current, the sum of the electron current and the hole current strictly obeys the knot law. To describe with an excellent approximation the behavior of a bipolar transistor, we can limit ourselves to the transport of only one type of charge: the minority carriers in the base. The node equation for these carriers must take into account their storage as well as their disappearance by recombination. This equation, also improperly called Kirchhoff's equation, has been improved to describe transistor lasers.

## METHODS OF CONNECTING THE RESISTORS

Any part of an electrical circuit communicates with the rest of the circuit through a number of terminals. The simplest situation is where the circuit part is a dipole. If the dipole is passive (does not contain generators) consisting only of resistors, then it can be replaced by a single resistor, called an equivalent resistor, so that the rest of the circuit does not "feel" the replacement.
A resistor is equivalent to a group of resistors if when the same voltage is applied to the resistance as the equivalent group, an electric current flows with the same intensity.

The laws of connection of resistances apply strictly to ohmic conductors:

- connecting resistors in series: $R_{s}=R_{1}+R_{2}+R_{3}$
- connecting resistors in parallel: $\frac{1}{R_{p}}=\frac{1}{R_{1}}+\frac{1}{R_{2}}+\frac{1}{R_{3}}$


## CONNECTION OF RESISTORS IN SERIES



If two or more resistors are connected in series they are part of the same branch of an electrical network.

Grouping resistors in series achieves the same power.
Considering a group of three resistors connected in series on each resistor is a voltage U1; U2 and
U3 we can write:
$\mathrm{U}=\mathrm{U} 1+\mathrm{U} 2+\mathrm{U} 3$
Using resistors, Ohm's law is written:
$\mathrm{U}=\mathrm{IR} 1+\mathrm{IR} 2+\mathrm{IR} 3$
or $\mathrm{U}=\mathrm{I}(\mathrm{R} 1+\mathrm{R} 2+\mathrm{R} 3)$
We apply the same law for the equivalent circuit:
$\mathrm{U}=\mathrm{IRS}$
This results in the following relationship: $\mathrm{R}=\mathrm{R} 1+\mathrm{R} 2+\mathrm{R} 3$

In the general case, when connecting resistors in series, it is written:

$$
R_{s}=R_{1}+R_{2}+R_{3} \text { or : } R_{s}=\sum_{k=1}^{n} R_{k} .
$$

The equivalent resistance is always greater than any of the component resistances Rk.

## CONNECTING RESISTORS IN PARALLEL



Two or more resistors connected in parallel are connected between the same two nodes

The equivalent resistance when connected in parallel is subjected to the same voltage $U$ at the terminals
According to Kirchhoff's first law
$\mathrm{I}=\mathrm{I} 1+\mathrm{I} 2+\mathrm{I} 3$ and

$$
\begin{aligned}
& I=\frac{U}{R_{1}}+\frac{U}{R_{2}}+\frac{U}{R_{3}} \text { or } \\
& I=U\left(\frac{1}{R_{1}}+\frac{1}{R_{2}}+\frac{1}{R_{3}}\right) \\
& \text { and } I=\frac{U}{R_{p}}
\end{aligned}
$$

For the equivalent circuit $\frac{1}{R_{p}}=\frac{1}{R_{1}}+\frac{1}{R_{2}}+\frac{1}{R_{3}}$
so: $\frac{1}{R_{p}}=\sum_{k=1}^{n} \frac{1}{R_{k}}$ if we have n resistors connected in parallel
the equivalent resistance Rp of the parallel connection is always lower than any resistance Rk If we have 2 resistances connected in parallel, the equivalent resistance is calculated with the relation:
$R_{p}=\frac{R_{1} \cdot R_{2}}{R_{1} \cdot+R_{2}}$

## PROCEDURE FOR THE PRACTICAL WORK - STEPS

1. Read the values of the voltage $U$ on the voltmeter and the values of the resistors $R$
2. The intensity of the electric current $I_{\text {read }}$ is determined
3. The calculation is performed using the following formula: $I_{\text {calculated }}=\frac{U}{\mathrm{R}}$

## EXPERIMENTAL PART

VERIFICATION OF OHM'S LAW FOR A SIMPLE CIRCUIT


1. Read the value of the electric voltage $U$ on the voltmeter and the value of the resistance R

For a resistance of $\mathrm{R}=1000 \Omega$
$R=1000 \Omega$
$\mathrm{U}=10 \mathrm{~V}$

Enter the read values in the following table:

|  | $\mathrm{U}(\mathrm{V})$ | $\mathrm{R}(\Omega)$ | $\mathrm{I}_{\text {read }}(\mathrm{A})$ | $\mathrm{I}_{\text {calculated }}(\mathrm{A})$ |
| :--- | :--- | :--- | :--- | :--- |
| 1 |  |  |  |  |
| 2 |  |  |  |  |
| 3 |  |  |  |  |
| 4 |  |  |  |  |
| 5 |  |  |  |  |
| 6 |  |  |  |  |
| 7 |  |  |  |  |

## CONNECTION OF RESISTORS IN SERIES


2. Enter the values read in the following table:

|  | u <br> $(\mathrm{T})$ | R <br> $(\Omega)$ | $\mathrm{I}_{\text {took }}$ <br> $(\mathrm{A})$ | $\mathrm{R}_{\text {total }=\mathrm{U} / \mathrm{I}_{\text {mon }}}^{(\Omega)}$ | $\mathrm{R}_{\text {total calculated }}$ <br> $(\Omega)$ | $\mathrm{I}_{\text {calculated }}$ <br> $(\mathrm{A})$ |
| :--- | :--- | :--- | :--- | :--- | :--- | :--- |
| 1 |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |
| 4 |  |  |  |  |  |  |
| 5 |  |  |  |  |  |  |

## CONNECTING RESISTORS IN PARALLEL

$\mathrm{R} 1=\mathrm{R} 2=\mathrm{R} 3=\mathrm{R}=1000 \Omega$

## $\mathrm{U}=10 \mathrm{~V}$


2. Enter the data in the following table:

|  | u <br> $(\mathrm{T})$ | R <br> $(\Omega)$ | $\mathrm{I}_{\text {read }}$ <br> (A) | $\mathrm{R}_{\text {total }}=\mathrm{U} / \mathrm{I}_{\text {read }}$ <br> $(\Omega)$ | $\mathrm{R}_{\text {total calculated }}$ <br> $(\Omega)$ | $\mathrm{I}_{\text {calculated }}$ <br> $(\mathrm{A})$ |
| :--- | :--- | :--- | :--- | :--- | :--- | :--- |
| 1 |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |
| 4 |  |  |  |  |  |  |
| 5 |  |  |  |  |  |  |


$R=1000 \Omega$
$\mathrm{U}=10 \mathrm{~V}$
2. Enter the data in the following table:

|  | u <br> $(\mathrm{T})$ | R <br> $(\Omega)$ | $\mathrm{I}_{\text {took }}$ <br> $(\mathrm{A})$ | $\mathrm{R}_{\text {total }}=\mathrm{U} / \mathrm{I}_{\text {mon }}$ <br> $(\Omega)$ | $\mathrm{R}_{\text {total calculated }}$ <br> $(\Omega)$ | $\mathrm{I}_{\text {calculated }}$ <br> $(\mathrm{A})$ |
| :--- | :--- | :--- | :--- | :--- | :--- | :--- |
| 1 |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |
| 4 |  |  |  |  |  |  |
| 5 |  |  |  |  |  |  |

## LIGHTING MEASUREMENTS

## NOTIONS OF PHOTOMETRY

Considering that the radiation field of light radiation is an integral part of the radiation field of optical radiation, we can make the following observations:
a ) the photometric quantities must have their correspondent between the radiometric quantities
b ) the units used for measuring radiometric quantities are based on the fundamental units used for measuring radiometric quantities are based on the fundamental units used for measuring length, mass and time, to which is added the complementary unit for measuring the solid angle.
c ) the units we will use to measure the photometric quantities must have the same fundamental units as the radiometric quantities to which the fundamental unit called CANDELA must be added.
d ) the CIE notations lead to the following conclusions: the size definition, the size symbol and the unit of measurement of photometric quantities were recommended and accepted by the international commission for lighting.

LIGHT from the point of view of the illuminator's technique is the radiant energy (or radiation) evaluated according to its ability to produce the visual sensation in people.

The assessment of radiant energy as light is done with the help of the eye, the primary optical instrument compared to which all other optical instruments are only self-adjoint.

If different spectral distributions are involved, the basic norms used to evaluate the radiant flux as luminous flux are given by a set of sensitivity factors of the eye. In order to distinguish between the mechanisms of the normal eye depending on the values of the radiation flow, we will call the vision during a clear day daytime vision, and the vision during evening twilight view. Daytime vision corresponds to relatively high radiation fluxes and twilight vision corresponds to relatively low radiation fluxes.

## LIGHT SPECTRAL YIELD

The sensation of vision varies a lot depending on the frequency of the radiation or the wavelength of the radiation in the light spectral range. The eye as a standard observer (CIE) adapted for daytime vision presents the maximum sensitivity for $\lambda=555 \mathrm{~nm}$. The visibility of the eye decreases relatively quickly on either side of this value, reaching the value 0 for: $\lambda=380 \mathrm{~nm}$ and $\lambda=$ 780 nm . When it is normalized in relation to its maximum value corresponding to the value $\lambda=555$ nm the visibility corresponding to a given wavelength $\lambda$ - we thus define the luminous spectral efficiency $K_{\lambda}$ (or relative brightness). $K_{\lambda}$ is the standard visibility factor of the human eye adapted for daytime vision and with the trichromatic component $y(y \lambda)$ of the standard observer within the trichromatic system (CIE).

The values of the luminous spectral yield have been accepted both by the International Commission for Illumination and by the International Union of Pure and Applied Physics and by the Weights and Measures Committee. They were used as the basis for establishing different types of photometric standards.
LUMINOUS FLUX F (CIE) is defined as the temporal rate of variation of light through a given surface

$$
\Phi=K \int_{0}^{\lambda} K_{\lambda} \Phi_{e \lambda} d \lambda
$$

where $\Phi_{\text {is it }}$ - is the spectral radiant flux
The integration must be considered for the wavelength domains where the integral is nonzero. Beyond the visible range the integrand is zero because $K_{\lambda}$ is null. Choosing the value of the proportionality constant $\mathrm{k}=683$ and expressing the radiant flux $\Phi_{\mathrm{e} \lambda}$ in watts, the unit of measurement for luminous flux $\Phi$ is the lumen (lm)

The lumen is the luminous flux in the solid angle unit of 1 steradian created by a point light source with the intensity of 1 candela. The lumen is the luminous flux through the surface unit whose points are located at a distance of 1 m from a light source of 1 candela.

The LUMINOUS EFFICIENCY of a light source is defined as the ratio between the luminous flux $\Phi_{\lambda}$ and the radiant flux $\Phi_{\mathrm{e} \lambda}$. The old name of luminous efficiency is LIGHT EFFICIENCY FACTOR . The unit of measurement of luminous efficiency is $1 \mathrm{~m} / \mathrm{W}$. The luminous efficiency is not the same as the efficiency of a radiation source, which is also expressed in $1 \mathrm{~m} / \mathrm{W}$, but which refers to the supply power of the radiation source and not to the radiant flux of the source.

## THE QUANTITY OF LIGHT $Q=\int_{0}^{t} \Phi d t \quad$ (C.I.E)

The amount of light $\mathbf{Q}$ is given by the time integral of the luminous flux
If the luminous flux is a constant quantity over time, the amount of light is given by the product of the luminous flux and the operating time of the light source. The unit of measure for the amount of light is $\mathbf{l m} . \mathbf{h}$. In the case of light sources that have a short operating time, such as those used in photography, the unit of measurement is lm.s.

THE LUMINOUS INTENSITY $I=\frac{d \Phi}{d \Omega} \quad$ (C.I.E.)
The luminous intensity of a light source in a given direction is defined as the density of the solid angle of the luminous flux from that direction.

The luminous intensity is equal to the ratio between the luminous flux through an area element perpendicular to the direction of the luminous flux and the solid angle element that the area element subtends in relation to the light source, considered point-shaped. Since the solid angle must have a peak, the definition given for the luminous intensity applies only to point sources. In many cases, the dimensions of the light sources are negligible compared to the distance at which they are observed, so that it can be considered that the light comes from a point source.

The unit of measurement of luminous intensity $I$ is the fundamental unit of measurement in photometry and bears the name of candella -cd - (CIE)
LIGHTING $E=\frac{d \Phi}{d S} \quad$ (C.I.E)

The surface density of the luminous flux $\Phi$ through a surface $\mathbf{S}$, given is called illumination $\mathbf{E}$. When quantitative values are not applied instead of lighting, the amount of lighting or illuminance is used. The unit of measure (CIE) used to measure illumination is lux (lx) and is defined as the
illumination of a surface with an area of $\mathbf{1} \mathbf{m}^{\mathbf{2}}$ on which there is a uniform distribution of a luminous flux equal to $1 \mathbf{l m}$. If the unit of area is $1 \mathrm{~cm}^{2}$, the unit of illumination is (CIE) is the photo(ph)

BRIGHTNESS

$$
B=\frac{d I}{d S_{n}}=\frac{d^{2} \Phi}{d S_{n} d \Omega}=\frac{d \Phi}{d \Omega}
$$

or luminance or brightness is defined as light intensity corresponding to the surface area unit of a light source viewed from a given direction.

The term brightness usually used refers to the color attribute of the surfaces from which light usually comes to the eye. With the help of brightness, visual perceptions can be classified from unclear (dark) to very bright.

The internationally recognized unit of measurement of brightness is cd.m ${ }^{-2}$

The brightness of a surface in a given direction can also be expressed in
$1 \mathrm{~m} . \mathrm{m}^{-2}$ provided that the respective surface is illuminated uniformly and perfectly diffused. A perfectly diffused surface must satisfy exactly the cosine law for both emission and reflection, which is difficult to achieve from a practical point of view. The unit of measure creates ambiguities of similarity between brightness and luminous intensity, while brightness is only the surface density of luminous intensity.

The luminous flux Фcorresponds to the area unit of a perfectly diffused surface whose light intensity in the direction normal to the surface is equal to $\mathbf{I}$ candela will be equal to $\pi$.I lumens. The brightness measurement unit accepted by (CIE) is the nit : $1 \mathrm{nt}=1 \mathrm{~cd} / \mathrm{m}^{2}$. The stilt is also used as a tolerated unit : $1 \mathrm{sb}=1 \mathrm{~cd} / 1 \mathrm{~cm}^{2}$.

The lambert is also used as a unit of brightness measurement : $\mathbf{1 L}=\mathbf{1} / \pi . \mathbf{c d} / \mathbf{c m}^{2}$. Lambert is the uniform brightness of a perfectly diffused surface that emits or reflects light at a temporal rate equal to $1 \mathrm{~lm} / 1 \mathrm{~cm}^{2}$

Mediation must to take into account both the variation of brightness from one point to another as well as the variation of brightness with the observation angle. The lambert is a unit of measure for intense light sources such as the Sun.

For less intense light sources, the submultiple called the millilambert can be used: mL.

LIGHT DENSITY $U=\frac{d Q}{d V}$
The amount of light corresponding to the volume unit is called light density $u$ and is measured in (lm.h)/(m $\mathbf{m}^{3}$ ) or (lm.s)/m ${ }^{\mathbf{3}}$

Among the photometric coefficients, the most important and most frequently used are:
a) LIGHT TRANSMISSION $\mathbf{t}$ is defined as the fraction of incident light transmitted by a certain medium.
b) LIGHT TRANSMISSION $\mathbf{r}$ is defined as the fraction of the incident light reflected by a certain medium.
c) LIGHT ABSORBANCE $\alpha$ is defined as the fraction of the incident light absorbed by a certain medium.

## MEASUREMENT STANDARDS. THE FUNDAMENTAL UNIT OF PHOTOMETRY

The standards used in photometry are primary standards and secondary standards.
CANDELLA (LUMINOUS INTENSITY MEASUREMENT UNIT) is used as the primary standard in photometry, although even the luminous flux unit could be more appropriate.

The first standards of light intensity were candles. At the beginning of this century, the proposal was made to use the black body radiator as a preliminary standard, but due to experimental difficulties, it was only in 1931 that the first established standard of this kind was achieved. The temperature used was that of platinum solidification. After several national laboratories made this black body standard and it was found that the result agreed satisfactorily, it was agreed in 1948 that it should be universally accepted as the primary light standard.

At the 13th General Conference of Weights and Measures in 1967, the definition of the candle was adopted:
CANDELLA is the luminous intensity in the perpendicular direction of a surface of $1 / 600000 \mathrm{~m}^{2}$ of the black body at the solidification temperature of platinum under a pressure of $101325 \mathrm{~N} / \mathrm{m}^{2}$.

Between the international illuminance used in 1948 as a unit of measurement of luminous intensity and the candella there is a relationship:

1 international lighting $=1.005$ candelas
The definition of the candle according to the black body was kept between the years 1948 1979 and the corresponding standard proved to be much more stable than its predecessors. Using Planck's law and the thermodynamic solidification temperature of platinum, the spectral power distribution could be calculated.

Two things led to the redefinition of the candle, the first referring to difficulties in the practical realization of the integral black body standard and the second referring to the rapid increase in the application of radiometric techniques in photometry. The redefinition of the candle in relation to the watt offers the possibility for several laboratories to build their own standards.

Candella is the luminous intensity in a given direction of a source that emits monochromatic radiation with frequency $v=540.10^{12} \mathrm{~Hz}$ and has a radiant intensity of (1/683) W/sr.

The value $v=540.10^{12} \mathrm{~Hz}$ corresponds to a wavelength ( $=555 \mathrm{~nm}$ )
Thus, the candle represents the fundamental unit in SI

Luminous spectral efficiency

$$
K(\lambda)=\frac{F_{\lambda}}{\Phi_{e \lambda}}
$$

$$
[K]_{\text {S.I. }}=\frac{l m}{W}
$$

The maximum value for $K(\lambda)$ is obtained for $\lambda=555 \mathrm{~nm}$ and is
$K(\lambda)=683 \mathrm{~lm} / \mathrm{W}$. The luminous spectral yield $\mathrm{V}(\lambda)$ of the radiant flux
(photopic vision) is equal to the luminous efficiency corresponding to a certain given wavelength and the luminous efficiency equal to $683 \mathrm{~lm} / \mathrm{W}$ corresponds to a wavelength $\lambda=555 \mathrm{~nm}$.

So we can write: $V(\lambda)=\frac{K(\lambda)}{683}$
Any photometric quantity $\mathrm{X}_{\mathrm{f}}$ (corresponding to daytime vision) can be obtained from a radiometric quantity $\mathrm{X}_{\mathrm{e}}$ using the relation:

$$
X_{f}=\int_{380 n m}^{780 n m} X_{e, \kappa} \cdot V(\lambda) \cdot d \lambda
$$

## 2. Lighting measurements.

The following values are recommended for lighting measurements

| The location |  | illumination |
| :---: | :---: | :---: |
| OFFICE | Conference rooms, reception rooms | 200-700 |
|  | Church premises | 700-1500 |
| factories (Companies) | Entrance passages, goods packaging departments | 150-300 |
|  | Production department (production line) | 300-750 |
|  | Quality department (verification) | 750-1500 |
|  | Department of assembly of electronic components | 1500-3000 |
| Hotels | The waiting halls | 100-200 |
|  | Reception | 200-1000 |
| shops <br> Supermarkets | Corridors, stairs, | 150-200 |
|  | Table for sale | 750-1500 |
|  | Showcases | 1500-3000 |
| HOSPITALS | Salons for patients | 100-200 |
|  | Examination room | 200-750 |
|  | Operating room | 750-1500 |
|  | Room for emergency treatments | 750-1500 |
| schools, universities | amphitheatres | 100-300 |
|  | Classrooms | 300-750 |
|  | laboratory | 750-1500 |

## EXPERIMENTAL ASSEMBLY:

## The measurements are made using the digital luxmeter BF06 TROTEC with photodetector

Figure 1. Digital luxmeter BF06 TROTEC with fotodetector

## PROCEDURE FOR THE PRACTICAL WORK - STEPS

Using the digital luxmeter BF06 device, light intensity measurements will be performed.
The practical work is done using the Digital Luxmeter BF06 device and a lamp

|  | E(lx) | $\mathrm{d}(\mathrm{x})$ |
| :--- | :--- | :--- |
| 1 |  |  |
| 2 |  |  |
| 3 |  |  |
| 4 |  |  |
| 5 |  |  |
| 6 |  |  |
| 7 |  |  |
| 8 |  |  |
| 9 |  |  |

Light intensity determinations will be made depending on the distance from the light source, according to the following function: $\mathrm{E}=\mathrm{E}(\mathrm{x})$

## MEASUREMENT OF SOUND INTENSITY

## THEORETICAL INTRODUCTION

Acoustics is the science of sound, considered to be that part of physics that deals with the study of phenomena regarding the production, propagation, reception of sounds as well as their effects. In current usage, the term "sound" refers not only to the phenomenon in the air responsible for the sensation of hearing, but also to humans or animals. Disturbances with low frequencies (infrasounds) or high frequencies (ultrasounds) that are received by a human or animal auditory organ are also considered "sound"; we can talk about underwater sound, sound in solids or sound in structures.

The wide field of acoustics and the area of interest are attributed to multiple applications such as:
a) the existence of the ubiquitous nature of mechanical disturbances generated by natural causes and human activity;
b) the existence of the sensation of hearing, of the human vocal ability, of communication through sound together with the variety of psychological influences that sound has on the auditory organ.

Fields such as speech, music, sound recording and reproduction, telephony, audiology, architectural acoustics and noise control have a strong association with the auditory sensation.

Sound is a means of transmitting information, independent of our natural ability to hear, it is a significant factor especially in underwater acoustics. The physical effects of sound on substances and bodies with which it interacts represent other areas of interest and application.

Figure 1 shows the fields of acoustics and the disciplines with which it is associated. The first ring shows the traditional subdivisions of acoustics and the outer one presents the technical and artistic fields in which acoustics are applied.


Figure 1. Schematic representation of fields of acoustics

## Historic

The idea that sound is a wave phenomenon was born from observations about water waves. The notion of a wave means an oscillatory disturbance that moves far from the source and does not transport any amount of matter at long propagation distances. The wave interpretation was also compatible with Aristotle's (384-322 B.H.) statement that air motion is generated by a source that pushes the surrounding air forward in such a way that sound propagates unaltered over long distances until the air disturbance is extinguished. A pertinent experimental result with antecedents from the time of Pythagoras (550Î.H) shows that the air movement generated by a body (source) that oscillates with a single musical note is also an oscillatory movement and has the same frequency as the body (source); historically this is related to the development of laws for the natural frequencies of vibrating strings and for the physical interpretation of musical consonants. The main authors of these interpretations were the French philosopher Marin Mersenne (1588-1648) who was called the father of acoustics and Galileo Galilei (1564-1642) whose book "Mathematical Discourses Concerning Two New Sciences" (1638) contains the most credible statement, given until then, regarding frequency equivalence. Mersenne's description in the book "Harmonic universelle" (1636) of the first absolute determination of the frequency of an audible tone (at 84 Hz ) indicates that he had already demonstrated that the ratio of the absolute frequencies of two oscillating strings radiating in space (air) a tone musical and its octaves is $1: 2$; the ratio of the oscillation frequencies is also $1: 2$, which is
compatible with the hypothesis of equivalence between the source in the air and the air movement frequency.

The analogy with water waves was strengthened by the belief that the air motion associated with musical sounds is oscillatory and by the observation that sound propagates with a finite speed.

Another known problem was that sound "bends" at corners, suggesting diffraction, a phenomenon often observed in water waves. Also, the classic experiment of Robert Boyle (1640) on the radiation of sound with the help of an alarm clock inserted in a partially empty glass vessel, ensures that air is necessary for the production and propagation of sound.

A special point of view on the wave was held by Gassendi (contemporary with Marsenne and Galileo) who argued that sound is due to a "storm of atoms" emitted by the body that emits sound, the speed of sound is the speed of atoms and the frequency represents the number of atoms emitted in time unit.

The apparent conflict between the ray and wave theories played an important role in the history of the sister science, optics, but the theory of sound developed from the beginning as a wave theory. When ray concepts were used to explain acoustic phenomena, as Reynolds and Rayleigh did in the 19th century, they were seen as mathematical approximations of wave theory. The successful incorporation of geometric optics into wave theory demonstrated that approximate, viable models of the wave phenomenon could be experienced in terms of ray concepts. (This recognition strongly influenced twentieth-century developments in architectural acoustics, underwater acoustics, and noise control. )

The mathematical theory of sound propagation shown in Isaac Newton's (1642-1727) work entitled "Principia" (1686) included a mathematical interpretation of sound as pressure pulses transmitted by surrounding fluid particles. The diagrams in figure 2 show the divergence of the wavefronts after they pass through a slit. The mathematical analysis was limited to waves of constant frequency because it required a number of devices, circuits and approximations that could influence the frequency variability. A substantial advance in the development of a viable theory of sound propagation
to be based on mathematical and physical concepts was produced in the 18th century by Euler (1707 - 1783), Lagrange (1736-1813) and d'Alembert (1717-1783) when field theory began to receive a defined mathematical structure ; modern theories can be seen as improvements of the theory developed by Euler and his contemporaries.


Figure 2. Sound qualities

The sound qualities designate the types of subjective effects corresponding to the characteristics quantities of the incident wave.
The sound has the following qualities:

- the height corresponding to the frequency;
- the strength corresponding to the intensity;
- the timbre corresponding to the harmonic content of a compound sound.


### 2.1 Sound pitch

This quality of sound was studied for the first time by Galileo Galilei and represents the quality of the auditory sensation determined by frequency. Due to this quality, two or more sounds of the same loudness can be differentiated through a subjective scale so that it may appear to the auditory organ higher (acute) or deeper (grave), that is, from high sounds to low sounds; however, corresponding to an objective scale of frequencies, results in high (high) and low (low) frequencies, respectively. In this case, two sounds are considered to have equal pitches if they have the same frequency.

Due to the frequency dependence of the height, the three types of sounds could be defined:

- infrasound ( $1 \mathrm{~Hz}-16 \mathrm{~Hz}$ );
- sounds ( $16 \mathrm{~Hz}-16 \mathrm{KHz}$ ) and ultrasound ( $16 \mathrm{KHz}-200 \mathrm{MHz}$ );
- hypersound (> 200 MHz ).

Theoretically, this frequency dependence leads to the application of the expression of the frequency of an acoustic source and in the case of this sound quality, being of the form:

$$
\begin{equation*}
\mathrm{f}=\frac{\omega}{2 \pi}=\frac{1}{2 \pi} \sqrt{\frac{\mathrm{k}}{\mathrm{~m}}} \tag{1}
\end{equation*}
$$

where $\mathbf{m}$ is the mass of the oscillating source that stores kinetic energy and $\mathbf{k}$ is the elasticity constant of the connecting element that stores potential energy.
Thus the production of a higher sound leads to a higher frequency of the wave process.
A solid body (material) possesses a series of own oscillation frequencies, one of which we can consider as fundamental, which determines the pitch of the sound and corresponds to a pure tone.
The human auditory organ can appreciate if two sounds are of the same pitch even though they are produced by different instruments and intensities. The causal correspondence between frequency and height is stable, so that at the same height, that is, the same sound as sensation. This property of fidelity of the human auditory organ to frequency made possible the existence of the musical scale and even music.

Thus, the pitch of the sound reflects the sensitivity of the human auditory organ to variations in frequency; with pure tones there is a difference of $0.2 \%$ in the frequency range $400 \mathrm{~Hz} \square \square 4000 \mathrm{~Hz}$ which manifests itself in the form of a change in height. Outside the considered frequency range, the sensitivity decreases with the frequency variation. This frequency-dependent sensitivity implies certain restrictions on the calibration and frequency adjustment of radios, musical or electronic instruments. The reception of the sound height is done logarithmically, and the height difference according to the frequency is determined according to the frequency ratio. According to the pitch scale, this frequency between the pitches of two tones with different frequencies is called an interval , and in music, a musical interval .

Experimentally, the pitch of the sound is determined by different methods, among which we mention the recording method and the beating method .

Method of registration . The oscillogram of the sound whose height we want so is recorded
we determine and on the same tape the oscillogram of a pendulum that beats the second is also recorded. In this way, the number of periods per second of the researched sound can be determined on the graph.

The method of beating. The phenomenon of "beating" is known, explained in the "Oscillations" chapter. The beat frequency $\mathbf{f b}$ is equal to the difference in the frequencies of the two compared sounds, i.e. $\mathrm{fb}=\mathrm{f} 2$ - f 1
where the sound of frequency f 1 is known, that of frequency f 2 must be determined and the frequency of beats fb is determined auditorily.

### 2.2 Intensity (loudness) of the sound

The sound intensity at a point of the sound field is defined as the amount of sound energy that traverses the unit of flat surface located at that point and oriented perpendicular to the direction of sound propagation in the unit of time.

The expression of intensity $\mathbf{I}$ is:

$$
\mathrm{I}=2 \pi^{2} \mathrm{f}^{2} \mathrm{~A}^{2} \rho \mathrm{~V}
$$

for a unitary surface, being proportional to the square of the local frequency and amplitude that characterizes the oscillatory system (source) and to two quantities that characterize the propagation medium, namely the density $\square \square$ and the speed of sound $v$.

Sound intensity is an objective physical quantity having a physical dimension

$$
[\mathrm{I}]=\frac{[\mathrm{F}] \cdot[\mathrm{L}]}{[\mathrm{L}]^{2} \cdot[\mathrm{~T}]}
$$

and as a unit of measure:

$$
\mathrm{I}_{\mathrm{SI}}=\frac{\mathrm{W}}{\mathrm{~m}^{2}}
$$

that is, it represents a power relative to the surface.
To define the acoustic pressure, we consider a continuous elastic medium for sound propagation.
If the wave phenomenon in this environment is null (there is no sound) then the local pressure is given only by the static pressure pst ; if longitudinal acoustic (sound) waves propagate through the
elastic medium, then at a point of the medium the total pressure $\mathbf{p t}$ will oscillate harmonically and the wave phenomenon appears. In this case, the dynamic pressure pdin , which would be due only to the wave effect, has the following expression: pdin $=\mathrm{pt}-\mathrm{pst}$ and depending on the propagation distance x and the wavelength $\square \square$ of the sound, the relation is obtained:

$$
\begin{equation*}
\mathrm{p}_{\mathrm{din}}=\mathrm{p}_{\max } \cos 2 \pi\left(\frac{\mathrm{t}}{\mathrm{~T}}-\frac{\mathrm{x}}{\lambda}\right) \tag{3}
\end{equation*}
$$

for a period T of the oscillations at a given time t and where p

$$
\mathrm{p}_{\max }=\frac{2 \pi \mathrm{AE}_{\mathrm{ac}}}{\lambda}
$$

The above expression (3) represents the momentary acoustic pressure at a certain distance x from the source of longitudinal acoustic waves and which acts on the auditory organ (tympanum) giving rise to the auditory sensation .

Depending on the maximum acoustic pressure (pmax), the acoustic intensity has the expression:

$$
\mathrm{I}=\frac{\mathrm{p}_{\max }^{2}}{2 \rho \mathrm{v}}
$$

(4)
and depending on the maximum speed of the material particles that make up the elastic medium vmax, there is the relationship:

$$
\begin{equation*}
\mathrm{I}=\frac{1}{2} \mathrm{v}_{\max }^{2} \rho \mathrm{v} \tag{5}
\end{equation*}
$$

These relations are important in practice, because they allow the determination of acoustic intensity based on measurements of acoustic pressure or the speed of material particles. Making the report

$$
\frac{\mathrm{p}_{\max }}{\mathrm{v}_{\max }}=\rho \mathrm{v}
$$

a law is obtained acoustically analogous to Ohm's law in electricity, if it is considered that the maximum acoustic pressure corresponds to the electric voltage, and the speed corresponds to the current intensity; therefore, an impedance, but acoustic, of the form appears:
$\mathrm{Z}=\square \mathrm{v}$ (7)
which represents the specific acoustic impedance of the elastic medium of propagation.

### 2.3 Sound timbre

Timbre has different definitions, but they all refer to the spectral content of the composite sound. Thus, according to STAS 1957-66, timbre is a characteristic of the auditory sensation that allows to distinguish different compound sounds having the same fundamental frequency and intensity but different spectral compositions and according to STAS 1957-74, timbre represents the quality of compound sounds of sensibly equal intensities to could be differentiated according to their spectral composition. From these definitions it follows that the timbre is a unique, structural characteristic of a sound that allows the identification and tracking of a certain source from a group of sound sources. The objective explanation of this sound quality, the timbre, can be achieved by studying the classification
sounds according to the result of their harmonic analysis or indirectly according to the auditory sensation it produces.

Thus, it can be defined as follows:
a) the pure sound (tone) characterized by a single frequency and amplitude and which corresponds the oscillation of the simplest material body (source). In the elongation-time graphic representation, the pure sound is a sinusoid and in the intensity-frequency representation, the pure sound is a vertical line drawn towards the respective frequency (Fig. 3)


Figure 3

In nature, pure sounds can be obtained very rarely; most often they can be obtained in the laboratory with
the help of tone generators and used in acoustic measurements and audiometry;
b) the musical sound - is produced by the periodic oscillation of a certain body. This sound is composed of a pure, fundamental sound followed by a certain number of harmonics that represent a multiple of the fundamental sound. So, the musical sound is a complex sound consisting of a fundamental and several harmonics obtained from a material body of a certain geometric shape (musical instruments). Thanks to the timbre, a musical note produced by a musical instrument can be distinguished from the same note produced by another instrument (violin and flute).
c) noise - it is an excellently complex sound, made up of very numerous but chaotically distributed components on the frequency band. The oscillogram of a noise does not present periodicity and its components are non-harmonic. As a physiological effect, noise produces an unpleasant auditory sensation.

In conclusion, all these three qualities of sound, pitch, intensity, timbre are defined only in the presence of the human auditory organ. Experimental studies have highlighted the fact that there are also secondary influences between loudness and frequency, between pitch and intensity, and timbre is influenced by intensity and frequency.

## 3 Introduction to physiological acoustics

The production and propagation of elastic waves from a source together with their reception in the form of sounds represents a continuous, linked, indissoluble process. All this received information is taken in real time by the human brain which also creates the sensation of sound. Thus, in order for an elastic wave (objective quantity) to be transformed into the sensation of sound (subjective quantity), frequency and sensitivity conditions must be met.

### 3.1 The human auditory organ

The human auditory organ (the ear) is a specialized organ that analyzes the excitations, sorts and classifies them according to frequency, giving the sensation of height, according to intensity, giving the sensation of loudness and according to the spectral composition giving the sensation of timbre. That is why the human ear is also a spectroscope or selective sensory analyzer because through a process of analysis and synthesis of external stimuli (elastic waves) it creates the sensation of sound.

Simply, the ear is an intermediate auditory organ that transforms elastic waves into the sensation of sound. Figure 4 explains the fact that the excitation (objective element) represented by the sound intensity Is, which corresponds to a certain sound pressure, is transforms through the ear, U , into a subjective element (perception) represented by the auditory intensity at.


Figure 4. Schematic representation between excitation and perception

As a psycho-physical system for sound reception and processing, the human auditory organ, the ear, has three main regions:

- external ear - Ue,
- middle ear - Um,
- inner ear - Ui
regions that are the seat of physical processes that justify their shape and dimensions.


## AUDIBILITY LIMITS

The human ear can perceive sounds if the elastic waves fall within certain limits of frequency, intensity and duration .

Frequency limits . The human auditory organ has an audibility range of $16 \mathrm{~Hz}(20 \mathrm{~Hz})-16 \mathrm{KHz}$ $(20 \mathrm{KHz})$ because it largely depends on the individual; with age the upper limit begins to decrease. Intensity limits . These depend mainly on the frequency but also on the individual. for a certain frequency there is a lower limit of sound intensity, but also an upper limit when the sounds produce a painful and harmful sensation. The highest sensitivity of the human ear is for normal sounds $(1000 \mathrm{~Hz}$ frequency sounds), because the difference between the upper and lower limits is maximum. For these normal sounds, the amplitude of the acoustic pressure on the eardrum is $\mathrm{p} \min =\mathbf{2 . 8 4 . 1 0}$ -
$5 \mathrm{~N}^{\prime} \mathrm{m} 2$ and the minimum sound intensity Ismin $=10^{-12} \mathrm{~W} / \mathrm{m} 2$. In air, the amplitude of air oscillations for a plane wave is $A \min =0.123 \AA$ for the frequency of 1000 Hz and an air density of $1.29 \mathrm{Kg} / \mathrm{m}^{3}$.

This lower limit Ismin $=10-12 \mathrm{~W} / \mathrm{m} 2$ represents the threshold of audibility which is the minimum value of acoustic pressure or sound intensity that the human ear still receives in the form of sound. Depending on the frequency, the audibility threshold has the minimum value (maximum sensitivity) in the range $1000 \mathrm{~Hz} \square \square 3000 \mathrm{~Hz}$, which is practically a constant value. Recent measurements have shown that the sensitivity of the human ear is maximum around the frequency of 2500 Hz and at neighboring frequencies the audibility threshold increases considerably.
Statistically, we can show that at $\mathrm{f}=50 \mathrm{~Hz}, \mathrm{PA}$ (audibility threshold) increases by $5.10^{5}$ times and at $\mathrm{f}=20 \mathrm{~Hz}$ PA increases 108 times compared to the minimum value. Also, the sensitivity of the human ear decreases 104 times compared to the maximum value, for $\mathrm{f}=100 \mathrm{~Hz}$ and at $\mathrm{f}=400 \mathrm{~Hz}$ it decreases only 10 times.
The upper threshold of audibility is defined as the maximum value of the pressure or
Ismax sound intensities corresponding to normal sound ( $\mathrm{f}=1000 \mathrm{~Hz}$ );
that is, pmax $=2.10^{2} \mathrm{~N} / \mathrm{m} 2$ again
Ismax $=10^{2} \mathrm{~W} / \mathrm{m} 2$.
At higher values, the person feels a strong and unpleasant pressure on the eardrum that turns into pain (pain threshold). The ratio between the sound intensity values for the two upper and lower thresholds is:

$$
\frac{\mathrm{Is}_{\max }}{\mathrm{Is}_{\min }}=\frac{10^{2}}{10^{-12}}=10^{14}
$$

and that of the pressures is:

$$
\frac{\mathrm{p}_{\max }}{\mathrm{p}_{\min }}=\frac{2 \cdot 10^{2}}{2 \cdot 10^{-5}}=10^{7}
$$

which shows that the field of auditory intensity is very wide. The curves that delimit the two thresholds of audibility (upper and lower) depending on the frequency are represented in figure 5 , which delimits the surface (domain) of audibility inside them. The size of this surface (domain) is an indicator of a person's hearing capacity (the surface decreases in people with hearing defects).


Figure 5. Schematic representation of the audibility domain

## Duration limits

For the sound to be heard, the oscillation that produces the elastic waves must have a certain minimum duration. This is estimated at 60 ms , which means for a normal sound ( $\mathrm{f}=1000 \mathrm{~Hz}$ ) a persistence on the eardrum of 60 periods. The human ear can perceive differences in duration between two sounds of up to 10 ms . If the incident elastic waves have a shorter duration then this one it is interpreted by the ear as a click with undefined characteristics.

## WEBER - FECHNER LAW

This law represents a link between sound intensity Is (objective element) and auditory intensity Ia (subjective element). Weber and Fechner showed that the variation of hearing intensity Ia is proportional to the relative variation of sound intensity Is and has the relation:

$$
\begin{equation*}
\mathrm{dIa}=\mathrm{k} \frac{\mathrm{dIs}}{\mathrm{Is}} \tag{8}
\end{equation*}
$$

Integrating the relation above, we get: $\mathrm{Ia}=\mathrm{k} \ln \mathrm{Is}+\mathrm{C}$
To determine the value of the constant C , write the relation for the sound intensity corresponding to the lower threshold of audibility (Ismin), i.e.: Ia0 $=\mathrm{k} \ln$ Ismin +C
but, as Ia $0=0$, the value of the integration constant $C$ is found as: $\mathrm{C}=-\mathrm{k} \ln$ Ismin By entering the value of the constant $C$ in the initial relationship, we obtain:

$$
\begin{equation*}
\mathrm{Ia}=\mathrm{k} \ln \mathrm{Is}-\mathrm{k} \ln \mathrm{Is} \mathrm{~s}_{\min }=\mathrm{k} \ln \frac{\mathrm{Is}}{\mathrm{Is}_{\min }} \tag{9}
\end{equation*}
$$

which represents the Weber - Fechner law, which is a logarithmic law specific to the sound intensity measurement system.

The Weber - Fechner law shows that if the sound intensity increases in geometric progression, then the auditory intensity increases in arithmetic progression.

## ACOUSTIC LEVELS

## Sound intensity level

Experiments in electroacoustics led to the establishment of a quantitative link between the intensity of objective excitation and the intensity of subjective excitation. Thus, the sound intensity level Ns was defined as the sound intensity related to the lower threshold of audibility (Ismin). But, as the expression of this ratio presents very large numbers, up to $10{ }^{14}$, it was agreed that the linear scale should be replaced by a logarithmic scale, so that the relation for the sound intensity level is:

$$
\begin{equation*}
\mathrm{Ns}=\log \frac{\mathrm{Is}}{\mathrm{Is}_{\min }} \tag{10}
\end{equation*}
$$

where we recall that Ismin represents the minimum sound intensity (the lower threshold of audibility) and has the value of Ismin $=10-12 \mathrm{~W} / \mathrm{m} 2$ for a frequency of 1000 Hz . The unit of measure for the sound intensity level is the Bell (B), after Graham Bell, the inventor of the telephone. Since the Bel is considered to be a unit of measurement that is too large, another unit is used which is the tenth part of the Bel, i.e. deciBel (dB) which represents the minimum level difference that can be perceived by the ear. In this case, the sound intensity level has the expression:

$$
\begin{equation*}
\mathrm{Ns}=10 \log \frac{\mathrm{Is}}{\mathrm{Is}_{\min }} \approx 20 \log \frac{\mathrm{p}_{\max }}{\mathrm{p}_{\min }}(\mathrm{dB}) \tag{11}
\end{equation*}
$$

The level of sound intensity corresponding to the threshold of painful sensation is:

$$
\begin{equation*}
\mathrm{Ns}=10 \log \frac{\mathrm{Is}}{\max } \mathrm{Is}_{\min }=10 \log 10^{14}=140 \mathrm{~dB} \tag{12}
\end{equation*}
$$

## The level of auditory intensity

Since the intensity of the auditory sensation depends on the sound intensity and the frequency of the sound, the auditory intensity level Na of a sound is defined as the sound intensity level of a normal sound ( $\mathrm{f}=1000 \mathrm{~Hz}$ ) which produces the same intensity of auditory sensation as the sound considered. This level of hearing intensity has the relationship:

$$
\begin{equation*}
\mathrm{Na}=10 \log \frac{\mathrm{Isn}}{\mathrm{Is}_{\min }} \approx 20 \log \frac{\left(\mathrm{p}_{\max }\right)_{\mathrm{n}}}{\mathrm{p}_{\min }} \tag{13}
\end{equation*}
$$

where Isn and (pmax)n represent the sound intensity and respectively the sound pressure amplitude of a normal sound of 1000 Hz . The unit of measure is the phon. A phon roughly corresponds to the resolving power of the ear relative to the loudness of the sound. The loudness of the sound expressed in phons is equal to the sound intensity level expressed in decibels of the normal sound of 1000 Hz which produces the same intensity of the auditory sensation.


Figure: the passing bands of the human ear and of some animal species.

## EXPERIMENTAL ASSEMBLY

A sound intensity measuring instrument is used in the montage, which records the sound intensity from second to second

## PROCEDURE FOR THE PRACTICAL WORK

Sound intensity determinations will be recorded and the data will be processed Sound intensity (dB) as a function of time (s)

The maps obtained for each registration will be presented

| No. crt. | $\mathrm{I}(\mathrm{dB})$ | $\mathrm{T}(\mathrm{s})$ |
| :--- | :--- | :--- |
| 1 |  |  |
| 2 |  |  |
| 3 |  |  |
| 4 |  |  |
| 5 |  |  |
| 6 |  |  |

The values of the maximum and minimum intensities and the average values recorded will be determined.

Data will be taken every 10 seconds with the device and graphed: $\mathrm{I}=\mathrm{I}(\mathrm{t})$
With the digital device for recording the intensity of the sound, automatic recordings are made from second to second. The data are processed in Excel and the graphs $\mathrm{I}=\mathrm{I}(\mathrm{t})$ are obtained.

## DETERMINATION OF THE REFRACTIVE INDEX OF TRANSPARENT SOLID AND LIQUID MEDIA USING THE MICROSCOPE

## INTRODUCTION

From the point of view of operation, the microscope is like a magnifying glass in the sense that the final image of an object is virtual and magnified, the only significant difference between the operation of a microscope and a magnifying glass is that in the case of the microscope the image is inverted from the real orientation of the object (as opposed to the magnifying glass).

Figure 1 represents the L lens in the position of the objective lens closest to the object. In order for the image of the object (source $S_{1}$ ) created by the microscope to be clear in the field of view of the microscope, the object (source $S_{1}$ ) must be located at a distance $\mathrm{x}_{1}$ determined by the characteristics of the objective lens. If a plane-parallel plate (medium with refractive index "n" greater than that of the medium) of thickness " d " is interposed between the object and the objective lens, the clarity of the image is compromised. To obtain a clear image again, the distance between the objective and the objective lens must be changed (reduced by " $\Delta x$ "), that is, the object from point $S_{1}$ must be brought to point $S_{2}$. The displacement " $\Delta x$ " imposed on the object, to restore the clear image, depends on the thickness " d " of the plane-parallel plate and its refractive index " n ". By determining the displacement " $\square \mathrm{x}$ " and having the thickness " d ", the refractive index of the plate can be calculated. The method is known as the " Chaulnes method" (Michel-Ferdinand d'Albert d'Ailly, duc de Chaulnes, 1712 - 1777; the method was proposed and described in 1768).


Figure 1. The influence of a plane-parallel plate positioned between the object and the objective lens


Figure 2. The apparent displacement of the image of an object due to a layer of refracting medium

Figure 2 illustrates the relationship between the object-lens distance ("d") for a clear image in the absence and presence of a homogeneous medium, of plane-parallel shape, thickness " d " and refractive index " n ". If in the absence of the plate the image of an object located at the point " x " " is clear, then, in the presence of the plane-parallel plate, the object must be moved to the position " x 2 ", with " $\Delta \mathrm{x}$ " closer to the lens, to obtain again a clear picture. With the help of the micrometer screw of the microscope, the " $x$ " displacement can be measured $\Delta$.

The length of the segment " AB " can be expressed from the right triangle " $\mathrm{ABS}_{1}$ ": $\mathrm{AB}=\mathrm{d} \bullet \mathrm{tg} i$ (1)
and from the right triangle " $\mathrm{ABS}_{2}$ " $: \mathrm{AB}=(\mathrm{d}-\Delta \mathrm{x}) \cdot \operatorname{tg} r$
By dividing, member by member, the two relations, it results:

$$
\frac{\operatorname{tg} r}{\operatorname{tg} i}=\frac{\mathrm{d}}{\mathrm{~d}-\Delta \mathrm{x}} \text { (3) }
$$

If the observation of the emergent ray is performed close to the S1S2A direction, the angles " $i$ " and " $r$ " are small; in this case the approximation is accepted: $\operatorname{tg} r \approx \sin r$ and $\operatorname{tg} i \approx \sin i$.

$$
\begin{equation*}
\frac{\operatorname{tg} r}{\operatorname{tg} i}=\frac{\mathrm{d}}{\mathrm{~d}-\Delta \mathrm{x}} \approx \frac{\sin r}{\sin i}=\mathrm{n} \tag{4}
\end{equation*}
$$

The refractive index " n " is expressed from relation (4):

$$
\mathrm{n}=\frac{\mathrm{d}}{\mathrm{~d}-\Delta \mathrm{x}}(5)
$$

## PROCEDURE FOR THE PRACTICAL WORK - STEPS

## I. CALIBRATION OF THE MICROMETRIC SCREW OF THE MICROSCOPE

To measure the change in the distance (" D ") between the objective lens and the sample table of the microscope, the micrometer screw located on the right side of the microscope is actuated. The wheel of the micrometric screw is divided into 200 parts (Figure 3)


Figure 3. Schematic representation of the disc divided into 200 parts of the micrometric screw

In order to be able to correlate the change in the distance " $D$ " with the measure of the rotation of the split disc, a body of known thickness ( 10 mm ) fitted with a thin wire "MN" is placed on the microscope table, Figure 4. The first time the body is placed on the microscope table like this so that the "MN" wire is in the upper part of the body. Rotate the divided disk until a clear image of the object is obtained in the visual field. The position of the split disc is noted.

The object in Figure 4 is inverted so that the "MN" thread ends up at the bottom. Thus, the objectlens distance was changed with the thickness of the object ( 10 mm ). Rotate the split disc of the micrometer screw to obtain a clear image again. The number of divisions between the two positions of the divided disk is noted. If between the two observations the disk has rotated with " N " complete rotations and with another " n " divisions, then $(200 \cdot \mathrm{~N}+\mathrm{n})$ divisions correspond to the change of 10 mm in the distance between the sample table and the objective. It means that a division on the divided

disk corresponds to a displacement " $\delta$ " expressed in millimeters by relation (6)

$$
\delta=\frac{10}{200 \cdot \mathrm{~N}+\mathrm{n}} \quad(\mathrm{~mm})(6)
$$

Figure 4. Disk for calibrating the micrometer screw of the microscope

## II. DETERMINATION OF THE REFRACTIVE INDEX OF A PLANE-PARALLEL GLASS PLATE

For reasons of convenience regarding the readings of the divisions of the divided disk of the micrometric screw, the operations are performed in the following sequence:

1) Measure the thickness of the glass plate, " $d$ ". On one of the faces of the plate is marked an "S $S_{1}$ " sign. Place the studied glass plate on the microscope table with the " $\mathrm{S}_{1}$ " sign on the lower surface. Actuate the micrometric screw until the clear virtual image " $\mathrm{S}_{2}$ " of the sign " $\mathrm{S}_{1}$ " is obtained (Figure 5-a). In this situation, the real distance between the sign " $\mathrm{S}_{1}$ " and the objective of the microscope is "x $x_{1}$ ". The distance between the apparent image " $\mathrm{S}_{2}$ " and the objective of the microscope is " $\mathrm{x}_{2}$ ".
2) Turn the glass plate so that the " $\mathrm{S}_{1}$ " mark reaches the upper surface of the glass plate. The real distance between the sign " $\mathrm{S}_{1}$ " and the objective of the microscope becomes " $\mathrm{x}_{1}$ ". The image of the sign " $\mathrm{S}_{1}$ " is now blurred (Figure 5-b).
3) Turn the micrometer screw to increase the distance between the " $\mathrm{S}_{1}$ " mark and the microscope objective. For the real distance " $\mathrm{x}_{2}$ " between the sign " $\mathrm{S}_{1}$ " and the objective of the microscope, the clear image of the sign "S 1 " is obtained again (Figure 5-c).
4) Read the difference between the divisions of the disc divided between configurations 1) and 2) and applying relation (6) calculate the displacement "D" of the sample table of the microscope.

It can be seen from Figure 5-c that the displacement "D" of the microscope table is
$\mathrm{D}=\mathrm{d}-\Delta \mathrm{x}$; from here: $\Delta \mathrm{x}=\mathrm{d}-\mathrm{D}$ (7)
5) Refractive index " n " is calculated based on relation (5):

$$
\mathrm{n}=\frac{\mathrm{d}}{\mathrm{~d}-\Delta \mathrm{x}}=\frac{\mathrm{d}}{\mathrm{~d}-(\mathrm{d}-\mathrm{D})}=\frac{\mathrm{d}}{\mathrm{D}}(8)
$$



Figure 5. Sequence of operations for determining the refractive index of a plane-parallel glass plate

## III. DETERMINATION OF THE REFRACTIVE INDEX OF A LIQUID USING THE MICROSCOPE

The refractive index of liquids can also be determined with the Chaulnes method. For this purpose, a glass vessel with internal thickness $d=1 \mathrm{~cm}$ can be used, a known accessory for spectrophotometric measurements (Figure 6-a). The tank can be closed tightly with the help of a plug made of Teflon.

The stages of determination are as follows:

1) Fill the vessel with the studied liquid (having the refractive index " $n$ ") and close it with the Teflon plug (Figure 6-b). It is important that the liquid enclosed in the tank does not contain air bubbles.
2) Create a mark on one of the outer faces of the cuvette and place the cuvette, in the lying position, on the sample stage of the microscope (Figure 6-c).
clear virtual image $S_{2}$ of the sign $S_{1}$ appears in the visual field of the observer. Note the position of the graduated disc of the micrometer screw.
3) Empty the tub, wash and dry. The empty cuvette is placed on the stage of the microscope. In the field of view of the observer the image is blurred.
4) Turn the graduated disc of the micrometric screw, in the direction indicated in Figure 3, until the image in the observer's visual field becomes clear (Figure 7-b). Note the complete number of rotations " D " and the number of divisions " y " with which the graduated disc was rotated.
5) The displacement " $\Delta x$ " is, this time, equal to the displacement " $D$ " of the sample table:
$\Delta \mathrm{x}=\mathrm{D}=200 \cdot \mathrm{~K}+\mathrm{y}(9)$
To calculate the refractive index " n " of the liquid in the tank due to the displacement " $\Delta \mathrm{x}$ ", a more general analysis of the system from an optical point of view is required.


Figure 6. Quartz vessel used for determining the refractive index of liquids


Figure 7. Determination of the refractive index of a liquid

Figure 8 represents the path of a light ray starting from the source "S" (the real object) and crossing three separation surfaces (the surfaces of the segments "AB", "CD" and "EF"). When crossing each surface, the limina undergoes refraction. The system models the two glass walls of the tank with the thickness " $\mathrm{d}_{1}$ " and " $\mathrm{d}_{2}$ " as well as with the refractive index " $\mathrm{n}_{1}$ " and " $\mathrm{n}_{2}$ ". Also, Figure 8 models the environment in the space between the two glass walls (thickness "d" and refractive index "n"). Refractive index of air " $\mathrm{n}_{0}$ " is considered equal to 1 .

Following the crossing of the layers and the multiple refraction of light, the observer (microscope lens) sees the virtual image " $\mathrm{S}_{2}$ " of the real object " $\mathrm{S}_{1}$ " displaced by the distance " $\Delta \mathrm{x}$ ".

Next, we aim to express the distance " $\Delta x$ " according to the thicknesses " $d_{1} ", " d ", d_{2}$ ", refractive indices " $\mathrm{n}_{1}$ ", " n ", " $\mathrm{n}_{2}$ ", " $\mathrm{n}_{0}$ " and the angles " $i_{1} ", " i_{2}$ ", " $r_{1} ", " r_{2}$ ".

From the triangle " $\mathrm{ABS}_{1}$ " the length of the segment " AB " is expressed, from the triangle " BCD " the length of the segment "CD" is expressed, and from the triangle "DEF" the length of the segment " EF " is expressed:
$\mathrm{AB}=\mathrm{d}_{1} \cdot \operatorname{tg} i_{1} \quad ; \quad \mathrm{CD}=\mathrm{d} \cdot \operatorname{tg} r_{1} \quad ; \quad \mathrm{EF}=\mathrm{d}_{2} \cdot \operatorname{tg} i_{2}(10)$
The distance " p " is the sum of the lengths of the segments " AB ", " CD " and " EF ": $\mathrm{p}=\mathrm{AB}+\mathrm{CD}+\mathrm{EF}(11)$

On the other hand, $\mathrm{p}=$ " $\mathrm{q} " \cdot \operatorname{tg} r_{2}$
From facts (11) and (11) it follows:
$\mathrm{p}=\mathrm{d}_{1} \cdot \operatorname{tg} i_{1}+\mathrm{d} \cdot \operatorname{tg} r_{1}+\mathrm{d}_{2} \cdot \operatorname{tg} i_{2}(13)$


Figure 8. Generation of the Virtual image " $\mathrm{S}_{2}$ " of the object " $\mathrm{S}_{1}$ " as a result of the consecutive refraction of light on three separating surfaces

If the observation direction is very close to the direction of the straight line " $\mathrm{S}_{1} \mathrm{~A}$ ", then all the angles " $i_{1}$ ", " $r_{1}$ ", " $i_{2}$ " and " $r_{2}$ " tend to zero and it can be assumed that:
$\mathrm{AB} \approx \mathrm{d}_{1} \cdot \sin i_{1} \quad ; \quad \mathrm{CD} \approx \mathrm{d} \cdot \sin r_{1} \quad ; \quad \mathrm{EF} \approx \mathrm{d}_{2} \cdot \sin i_{2}(14)$
Accepting this approximation, relations (12) and (13) become:
$\mathrm{p} \approx \mathrm{q} \cdot \sin r_{2}$ and $\mathrm{p}=\mathrm{d}_{1} \cdot \sin i_{1}+\mathrm{d} \cdot \sin r_{1}+\mathrm{d}_{2} \cdot \sin i_{2}(15)$
Applying in turn the Willebrord Snel van Roijen relation, related to refraction, to the surfaces "AB", "CD" and "EF", it results:

$$
\left\{\begin{array}{l}
\mathrm{n}_{1} \cdot \sin i_{l}=\mathrm{n} \cdot \sin r_{1} \\
\mathrm{n} \cdot \sin r_{l}=\mathrm{n}_{2} \cdot \sin i_{2} \\
\mathrm{n}_{2} \cdot \sin i_{2}=\mathrm{n}_{\mathrm{o}} \cdot \sin r_{2}
\end{array}\right\} \Rightarrow \mathrm{n}_{1} \cdot \sin i_{1}=\mathrm{n}_{\mathrm{o}} \cdot \sin r_{2}(16)
$$

From relations (6) it follows:

$$
\begin{equation*}
\sin i_{l}=\frac{\mathrm{n}_{\mathrm{o}}}{\mathrm{n}_{1}} \cdot \sin r_{2} ; \sin r_{l}=\frac{\mathrm{n}_{\mathrm{o}}}{\mathrm{n}} \cdot \sin r_{2} ; \sin i_{2}=\frac{\mathrm{n}_{\mathrm{o}}}{\mathrm{n}_{2}} \cdot \sin r_{2} \tag{17}
\end{equation*}
$$

Substituting relations (17) into relations (15), we obtain:

$$
\begin{equation*}
\mathrm{p} \approx \mathrm{q} \cdot \sin \mathrm{r}_{2} ; \mathrm{p} \approx \mathrm{n}_{\mathrm{o}} \cdot \sin r_{2} \cdot\left(\frac{\mathrm{~d}_{1}}{\mathrm{n}_{1}}+\frac{\mathrm{d}}{\mathrm{n}}+\frac{\mathrm{d}_{2}}{\mathrm{n}_{2}}\right) \tag{18}
\end{equation*}
$$

From the two relations (8) it follows:

$$
\begin{equation*}
\mathrm{p} \approx \mathrm{n}_{\mathrm{o}} \cdot\left(\frac{\mathrm{~d}_{1}}{\mathrm{n}_{1}}+\frac{\mathrm{d}}{\mathrm{n}}+\frac{\mathrm{d}_{2}}{\mathrm{n}_{2}}\right) \tag{19}
\end{equation*}
$$

Since $\square \mathrm{x}=\mathrm{d}_{1}+\mathrm{d}+\mathrm{d}_{2}-\mathrm{q}$, we obtain:

$$
\Delta \mathrm{x} \approx \mathrm{~d}_{1}+\mathrm{d}+\mathrm{d}_{2}-\mathrm{n}_{\mathrm{o}} \cdot\left(\frac{\mathrm{~d}_{1}}{\mathrm{n}_{1}}+\frac{\mathrm{d}}{\mathrm{n}}+\frac{\mathrm{d}_{2}}{\mathrm{n}_{2}}\right)(20)
$$

and after rearranging relation (2), the expression of displacement " $\square \mathrm{x}$ " results:

$$
\begin{equation*}
\Delta \mathrm{x} \approx \mathrm{~d}_{1} \cdot\left(1-\frac{\mathrm{n}_{\mathrm{o}}}{\mathrm{n}_{1}}\right)+\mathrm{d} \cdot\left(1-\frac{\mathrm{n}_{\mathrm{o}}}{\mathrm{n}}\right)+\mathrm{d}_{2} \cdot\left(1-\frac{\mathrm{n}_{\mathrm{o}}}{\mathrm{n}_{2}}\right) \tag{21}
\end{equation*}
$$

By specifying relation (21) for the conditions of the experiment:

$$
\mathrm{n}_{1}=\mathrm{n}_{2}=\mathrm{n}_{\text {bottle }} ; \mathrm{d}_{1}=\mathrm{d}_{2}=\mathrm{d}_{\text {glass }} ; \mathrm{n}_{\mathrm{o}}=\mathrm{n}_{\text {air }}=1
$$

the general relation (21) is simplified.
For the case where the vessel is filled with the liquid sample (refractive index " n "), relation (21) reduces to form (22).
$\Delta x_{\text {filled }} \approx 2 \cdot d_{\text {glass }} \cdot\left(1-\frac{1}{n_{\text {glass }}}\right)+d \cdot\left(1-\frac{1}{n}\right)$
For the case when the cuvette is empty ( $\mathrm{n}_{\text {air }}=1$ ), it follows:

$$
\begin{equation*}
\Delta x_{\text {empty }} \approx 2 \cdot d_{\text {glass }} \cdot\left(1-\frac{1}{n_{\text {glass }}}\right) \tag{2.3}
\end{equation*}
$$

Changing the position of the virtual image " $\mathrm{S}_{2}$ " (" $\Delta \Delta \mathrm{x}$ ") for the measurement performed with a filled cuvette and an empty cuvette,

$$
\begin{equation*}
\Delta \Delta \mathrm{x}=\mathrm{d} \cdot\left(1-\frac{1}{\mathrm{n}}\right) \tag{24}
\end{equation*}
$$

where " d " is the thickness of the space between the two glass walls (for the cuvette used in the experiment $\mathrm{d}=1 \mathrm{~cm}$ ). Explaining from relation (24) the refractive index of the liquid contained in the cuvette, we obtain:

$$
\begin{equation*}
\mathrm{n}=\frac{\mathrm{d}}{\mathrm{~d}-\Delta \Delta \mathrm{x}} \tag{25}
\end{equation*}
$$

that is, a relation of similar form to the first equality in expression (8). It is noted that it is not necessary to know the thickness of the glass walls ( $\mathrm{d}_{1}$ and $\mathrm{d}_{2}$ ).

If in the system represented in Figure 8, the thicknesses " d " and " $\mathrm{d}_{2}$ " tend to zero, then the situation created corresponds to the one represented in Figure 2. In this case, the general relationship (21) takes the particular form

$$
\begin{equation*}
\Delta \mathrm{x} \approx \mathrm{~d}_{1} \cdot\left(1-\frac{1}{\mathrm{n}_{1}}\right)_{\mathrm{g}} \tag{26}
\end{equation*}
$$

and from here the index of refraction of the glass plate (" $\mathrm{n}_{1}$ ") is expressed

$$
\begin{equation*}
\mathrm{n}_{1}=\frac{\mathrm{d}_{1}}{\mathrm{~d}_{1}-\Delta \mathrm{x}} \tag{27}
\end{equation*}
$$

through a relationship similar to (5).

## EXPERIMENTAL ASSEMBLY



Ocular (1) - contains the lenses ocular, typical magnification values: 5x, 10x and 2 x

Objective (2); typical values of magnifications:
4x, 5x, 10x, 20x, 40x, 50x and 100x; interchangeable

Table (3) - supports the test of examination
$\square$ Light source (4) - located under the table

Condenser (5) - system of lenses and diaphragms which controls the amount of light by focusing the light

Macrovisa (6) and microvisa (7) used to clarify image

## EXPERIMENTAL DATA PROCESSING


$x^{\prime}=x-d+d x$
$x^{\prime}=x+d(n-1)$
$x^{\prime}-x=\Delta x=d(n-1)$
$n=1+\frac{\Delta x}{d}$
$\mathrm{n}=1+\frac{\Delta \mathrm{x}}{\mathrm{d}}$
$\Delta x=200 K+y$
$K$ - the number of complete rotations again
$y$ - the number of divisions
$\mathrm{d}=1 \mathrm{~cm}$ - the side of the cuvette where the liquid is introduced

## 1. DETERMINATION OF THE REFRACTIVE INDEX IN TRANSPARENT SOLID MEDIUMS

## PROCEDURE FOR THE PRACTICAL WORK - STEPS

1. Calibrate the microscope
2. The conversion factor of the microscope can be determined by superimposing the image formed by the ocular microscope and the objective microscope.
3. A transparent plate with refractive index $n$ and thickness $d$ is placed on the microscope table. The value of $d$ is determined and entered in the table of values.
4. An image is created on the surface of the plate
5. Read the value $\Delta x=200 K+y$ from the microscope,
6. where K are the number of complete rotations and y are the number of divisions
7. The data are entered in the table
8. Calculations are made for the refractive index n with the following formula:
$\mathrm{n}=1+\frac{\Delta \mathrm{x}}{\mathrm{d}}$
$\Delta x=200 K+y$
$K$ - the number of complete rotations again

## $y$ - the number of divisions

$\mathrm{d}=\mathbf{1 0} \mathbf{~ m m}$
Table 1: determination of the refractive index $n$ for different blades of dimensions $d$

| No. | d (mm) | K (wheel) | Y (div) | $\Delta \mathbf{x}(\mu \mathbf{m})$ | $\mathbf{n}$ |
| :--- | :--- | :--- | :--- | :--- | :--- |
| 1 |  |  |  |  |  |
| 2 |  |  |  |  |  |
| 3 |  |  |  |  |  |
| 4 |  |  |  |  |  |
| 5 |  |  |  |  |  |

Plot $\mathrm{n}=\mathrm{n}(\mathrm{d})$ graphically

## 2. DETERMINATION OF THE REFRACTIVE INDEX IN TRANSPARENT LIQUID MEDIUMS

### 2.2 VARIATION OF REFRACTORY INDEX WITH CONCENTRATION PROCEDURE FOR THE PRACTICAL WORK - STEPS

1. Calibrate the microscope
2. The conversion factor of the microscope can be determined by superimposing the image formed by the ocular microscope and the objective microscope.
3. A transparent plate with refractive index $n$ and thickness $d$ is placed on the microscope table. The value of $\mathrm{d}=1 \mathrm{~cm}$ is known and entered in the table of values.
4. An image is created on the surface of the plate
5. Read the value $\Delta x=200 K+y$ from the microscope,
6. where K are the number of complete rotations and y are the number of divisions, $\mathrm{d}=1 \mathrm{~cm}$ - the side of the tank where the liquid is introduced
7. The data are entered in the table
8. Calculations are made for the refractive index n with the following formula:
$\mathrm{n}=1+\frac{\Delta \mathrm{x}}{\mathrm{d}}$
$\Delta x=200 K+y$
$K$ - the number of complete rotations again
$\mathbf{y}$ - the number of divisions
$\mathrm{d}=\mathbf{1 ~ c m}$ - the side of the tank where the liquid is introduced

TABLE 2: determination of the refractive index $n$ for the ethylene glycol solution of different concentrations

| $\mathbf{C}(\%)$ | K (wheel) | Y (div) | $\Delta \mathbf{x}(\mu \mathbf{m})$ | $\mathbf{n}$ |
| :--- | :--- | :--- | :--- | :--- |
| 100 |  |  |  |  |
| 90 |  |  |  |  |
| 80 |  |  |  |  |
| 70 |  |  |  |  |
| 60 |  |  |  |  |
| 50 |  |  |  |  |
| 40 |  |  |  |  |
| 30 |  |  |  |  |
| 20 |  |  |  |  |
| 10 |  |  |  |  |

Plot $\mathbf{n}=\mathbf{n}(\mathbf{c})$ graphically

### 2.3 VARIATION OF REFRACTIVE INDEX WITH DENSITY

 PROCEDURE FOR THE PRACTICAL WORK1. Calibrate the microscope
2. The conversion factor of the microscope can be determined by superimposing the image formed by the ocular microscope and the objective microscope.
3. A transparent plate with refractive index $n$ and thickness $d$ is placed on the microscope table. The value of $d$ is known and entered in the table of values.
4. An image is created on the surface of the plate
5. value is read $\Delta$ from the microscope.
6. The data are entered in the table
7. Calculations are made for the refractive index n with the following formula:
$\mathrm{n}=1+\frac{\Delta \mathrm{x}}{\mathrm{d}}$
$\Delta \mathrm{x}=200 \mathrm{~K}+\mathrm{y}$
$K$ - the number of complete rotations again
$\mathbf{y}$ - the number of divisions
$\mathrm{d}=1 \mathrm{~cm}$ - the side of the tank where the liquid is introduced

TABLE 2: Determination of the refractive index for solvents of different densities

| No. <br> crt | Substance | $\rho$ <br> $\left(\mathbf{k g} / \mathbf{m}^{3}\right)$ | K <br> $($ rotate $)$ | y <br> $($ div $)$ | $\Delta \mathbf{x}=\mathbf{2 0 0 ~ K + y}$ <br> $(\mu \mathbf{~ m})$ | $\mathbf{n}$ |
| :--- | :--- | :--- | :--- | :--- | :--- | :--- |
| 1 | Water $\left(\mathrm{H}_{2} \mathrm{O}\right)$ | 1000 |  |  |  |  |
| 2 | isopropyl alcohol | 786 |  |  |  |  |
| 3 | ethyl alcohol $\left(\mathrm{C}_{2} \mathrm{H}_{6} \mathrm{O}\right)$ | 789 |  |  |  |  |
| 4 | methanol $\left(\mathrm{CH}_{4} \mathrm{O}\right)$ | 1015 |  |  |  |  |
| 5 | aniline $\left(\mathrm{C}_{6} \mathrm{H}_{7} \mathrm{H}\right)$ | 1010 |  |  |  |  |
| 6 | acetic acid $\mathrm{C}_{2} \mathrm{H}_{4} \mathrm{O}_{2}$ | 1050 |  |  |  |  |
| 7 | DMSO | 1100 |  |  |  |  |
| 8 | ethylene $\operatorname{glycol} \mathrm{C}_{2} \mathrm{H}_{4} \mathrm{O}_{2}$ | 1110 |  |  |  |  |
| 9 | Glycerin | 1260 |  |  |  |  |

It is drawn graphically $n=n(\rho)$

## DETERMINATION OF THE FOCAL LENGTH OF LENSES USING THE BESSEL METHOD

## THEORETICAL INTRODUCTION

The optical behavior of thin lenses (where the thickness is much smaller than the diameter) can be characterized by the value of the focal length " $f$ ". If it is limited to the incident rays close to the main optical axis (paraxial rays), then it can be admitted that the focal length (the distance from the optical center of the lens to the point where the emerging ray intersects the main optical axis) does not depend on the distance between the emerging ray and the main optical axis of the lens.

The optical behavior of a lens is realized by its ability to create the image of an object. If the image is formed at the intersection of the real emergent rays, then the image is "real". If the image created by the lens is formed at the intersection of the extension of the real emergent rays, then the image is "virtual". Only a real image can be projected on a screen.

In the case of a real object, the distance at which the image is formed, as well as the nature of the image (real or virtual) depends on the focal length of the lens.

The location of the object and image relative to the lens is indicated by the "coordinates" of the object and image in relation to a reference system that has, as one of the coordinate axes ("x" axis), the main optical axis of the lens. The positive direction of the "x" axis is determined by the direction of the light ray starting from the real object. The origin of the "x" axis is chosen in the optical center of the lens. Figure 1 illustrates the sign of the " $\mathrm{x}_{\mathrm{o}}$ " and " $\mathrm{x}_{\mathrm{i}}$ " coordinates of the respective object and image.


Fig.
1

## Sign

convention
regarding
object and image coordinates in the case of a converging lens

The focal length of a converging lens is considered positive, while the focal length of a diverging lens is negative. In the example illustrated in Figure 1, the " $\mathrm{x}_{\mathrm{o}}$ " coordinate of the object is negative ( $\mathrm{x}_{\mathrm{o}}<0$ ), and that of the place where the real image is formed, " $\mathrm{x}_{\mathrm{i}}$ ", is positive $\left(\mathrm{x}_{\mathrm{i}}>0\right)$. If reference is made to the lens-object (" $\mathrm{L}_{\mathrm{o}}$ ") and lens-image (" $\mathrm{L}_{\mathrm{o}}$ ") "distances", these quantities are defined as positive, so $L_{o}=-x_{o}$ while $L_{i}=x_{i}$. (it is important to distinguish the notions of "coordinate" and "distance").

For a position of the lens, for which the clear image of an object (located at the " $\mathrm{x}_{\mathrm{o}}$ " coordinate) is formed on a screen (located at the "xi" coordinate), the relation of conjugate points (1) is respected.

$$
\frac{1}{\mathrm{f}}=\frac{1}{\mathrm{x}_{\mathrm{i}}}-\frac{1}{\mathrm{x}_{\mathrm{o}}}(1)
$$

If the distance between the object and the screen (where the image of the object is formed) is "D", this can be expressed by relation (2).

$$
\mathrm{D}=\mathrm{L}_{\mathrm{o}}+\mathrm{L}_{\mathrm{i}}=\left(-\mathrm{x}_{\mathrm{o}}\right)+\mathrm{x}_{\mathrm{i}}(3)
$$

" $x_{0}$ " and " $x_{i}$ " can be expressed as follows:

$$
\mathrm{x}_{\mathrm{i}} \cdot\left(-\mathrm{x}_{\mathrm{o}}\right)=\mathrm{D} \cdot \mathrm{f}(4)
$$

Relations (3) and (4) respectively express the sum and the product of two quantities ("- $x_{0}$ " and " $\mathrm{x}_{\mathrm{i}}$ ") depending on the distance " D " and the focal distance " f ". It is possible to write an equation of
the second degree in relation to the unknown "x", which has, as roots, precisely the quantities "- $\mathrm{x}_{0}$ " and " $\mathrm{x}_{\mathrm{i}}$ ":

$$
\begin{equation*}
x^{2}-D \cdot x+D \cdot f=0 \tag{5}
\end{equation*}
$$

If a certain relationship is observed between the quantities " D " and " f " (relation (8)), then the second degree equation has two distinct roots, " $\mathrm{x}_{\mathrm{A}}$ " and $" \mathrm{x}_{\mathrm{B}}$ ":

$$
\begin{align*}
& x_{A}=\frac{D+\sqrt{D^{2}-4 \cdot D \cdot f}}{2}  \tag{6}\\
& x_{B}=\frac{D-\sqrt{D^{2}-4 \cdot D \cdot f}}{2}
\end{align*}
$$

The condition for the roots (6) to be real and disjoint is for the discriminant to be positive, so $\sqrt{D^{2}-4 \cdot D \cdot f}>0$

The relation (7) is equivalent to the relation (8)

$$
\begin{equation*}
\mathrm{D}>4 \cdot \mathrm{f} \tag{8}
\end{equation*}
$$

which expresses the condition that there are two different positions of the lens for which a clear image of the object is formed at the distance " D " from the object.

Numerous methods are known for experimentally determining the focal length of a lens. An ingenious method is known as "Bessel 's method ". To understand the essence of this method, we will consider an optical assembly consisting of a thin converging lens, a real object and the real image of the object projected on a screen. If the object-screen distance is kept constant, two positions of the lens can be found, between the object and the screen, for which the clear image of the object is formed on the screen (Figure 2). In one of the positions of the lens, the enlarged image of the object is formed on the screen (Figure 2-A), and in the other position of the lens, the actual reduced image of the object is obtained (Figure 2-B). The coordinate of the object in the first case is " $\mathrm{x}_{01}$ " $\left(\mathrm{x}_{01}<0\right)$, and that of the image is " $\mathrm{x}_{\mathrm{i} 1}$ " $\left(\mathrm{x}_{\mathrm{i} 1}>0\right)$. The coordinate of the object in the second case is " $\mathrm{x}_{02}$ " $\left(\mathrm{x}_{02}<0\right)$, and that of the image is " $\mathrm{x}_{\mathrm{i} 2}$ " $\left(\mathrm{x}_{\mathrm{i} 2}>0\right)$. In the first case, the lens-object distance is $\mathrm{L}_{\mathrm{ol}}=-\mathrm{x}_{\mathrm{o}}$, and the lensimage distance is $L_{i 1}=x_{i 1}$. For the second position of the lens (Figure 2-B) we have $L_{o 2}=-x_{02}$ and $\mathrm{L}_{\mathrm{i} 2}=\mathrm{X}_{\mathrm{i} 2}$.


Fig. 2 Bessel's method for determining the focal length of a lens

The coordinates of the location of the object and the image for the two positions of the lens, at which the clear image of the object is obtained, are:

$$
\begin{align*}
-x_{o 1} & =\frac{D-\sqrt{D^{2}-4 \cdot D \cdot f}}{2} \quad ; \quad-x_{02}=\frac{D+\sqrt{D^{2}-4 \cdot D \cdot f}}{2} \\
x_{i 1} & =\frac{D+\sqrt{D^{2}-4 \cdot D \cdot f}}{2} \quad ; \quad x_{i 2}=\frac{D-\sqrt{D^{2}-4 \cdot D \cdot f}}{2} \tag{9}
\end{align*}
$$

The distance " d " between the positions of the lens for which the real and clear image of the object is obtained, is (Figure 2):

$$
\mathrm{d}=\mathrm{D}-\mathrm{L}_{\mathrm{o} 1}-\mathrm{L}_{\mathrm{i} 2}=\mathrm{D}-\left(-\mathrm{x}_{\mathrm{o} 1}\right)-\mathrm{x}_{\mathrm{i} 2}(10)
$$

Substituting in relation (1) the expressions found for the coordinates " $\mathrm{x}_{\text {o1 }} ", \mathrm{x}_{\text {i1 }}$ ", "x o2 " and "x i2 ", we obtain:
$\mathrm{d}=\mathrm{D}-\frac{\mathrm{D}-\sqrt{\mathrm{D}^{2}-4 \cdot \mathrm{D} \cdot \mathrm{f}}}{2}-\frac{\mathrm{D}-\sqrt{\mathrm{D}^{2}-4 \cdot \mathrm{D} \cdot \mathrm{f}}}{2}(11)$
from which it follows, for the distance " d ":

$$
d=\sqrt{D^{2}-4 \cdot D \cdot f}(12)
$$

The focal length of the lens can be explained from relation (12):

$$
\mathrm{f}=\frac{\mathrm{D}^{2}-\mathrm{d}^{2}}{4 \cdot \mathrm{D}}(13)
$$

Practically, the distance between the two positions of the lens for which the clear image of the object is formed is determined and the focal distance is calculated with relation (13).

## PROCEDURE

Figure 3 represents a setup for performing measurements.


Fig. 3 The optical assembly for determining the focal distance
Five determinations are performed, each with a different value of the distance "D" and - implicitly - resulting in other "d" values. For each pair of "D" and "d" values, the focal distance is calculated with relation (13). The average of the " $f$ " values and the standard deviation of the individual values around the average value are determined.

The data goes into a table according to the following model.

Table 1 Results

| no. | $\mathbf{D}(\mathbf{m m})$ | $\mathbf{L}_{\mathbf{o 1}(\mathbf{m m})}$ | $\mathbf{L}_{\mathbf{i 1}}(\mathbf{m m})$ | $\mathbf{L}_{\mathbf{o 2}(\mathbf{m m})}$ | $\mathbf{L}_{\mathbf{i 3}}(\mathbf{m m})$ | $\mathbf{d}(\mathbf{m m})$ | $\mathbf{f}(\mathbf{m m})$ |
| :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- |
| 1 |  |  |  |  |  |  |  |
| 2 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  |  |
| 4 |  |  |  |  |  |  |  |
| 5 |  |  |  |  |  |  |  |

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