ENERGY, ENVIRONMENT, BIOLOGY and BIOMEDICINE

Proceedings of the 2014 International Conference on Energy, Environment, Ecosystems and Development II (EEED '14)

Proceedings of the 2014 International Conference on Biology and Biomedicine II (BIO '14)

Prague, Czech Republic
April 2-4, 2014

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Keynote Lecture 1

Interpolation and Projective Representation in Computer Graphics, Visualization and Games

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Abstract: Today's engineering problem solutions are based mostly on computational packages. However, the computational power doubles in 18 months. In 15 years perspective the computational power will be of $2^{10} = 1024$ of today's computational power. Engineering problems solved will be more complicated, complex and will lead to a numerically ill conditioned problems especially in the perspective of today available floating point representation and formulation in the Euclidean space. Homogeneous coordinates and projective geometry are mostly connected with geometric transformations only. However, the projective extension of the Euclidean system allows reformulation of geometrical problems which can be easily solved. In many cases quite complicated formulae are becoming simple from the geometrical and computational point of view. In addition they lead to simple parallelization and to matrix-vector operations which are convenient for matrix-vector hardware architecture like GPU.

In this short tutorial we will introduce "practical theory" of the projective space and homogeneous coordinates. We will show that a solution of linear system of equations is equivalent to generalized cross product and how this influences basic geometrical algorithms. The projective formulation is also convenient for computation of barycentric coordinates, as it is actually one cross-product implemented as one clock instruction on GPU. Selected examples of engineering disasters caused by non-robust computations will be presented as well.

Brief Biography of the Speaker: Prof. Vaclav Skala is a Full professor of Computer Science at the University of West Bohemia, Plzen, Czech Republic. He received his Ing. (equivalent of MSc.) degree in 1975 from the Institute of Technology in Plzen and CSc. (equivalent of Ph.D.) degree from the Czech Technical University in Prague in 1981. In 1996 he became a full professor in Computer Science. He is the Head of the Center of Computer Graphics and Visualization at the University of West Bohemia in Plzen (http://Graphics.zcu.cz) since 1996.

Prof. Vaclav Skala is a member of editorial board of The Visual Computer (Springer), Computers and Graphics (Elsevier), Machine Graphics and Vision (Polish Academy of Sciences), The International Journal of Virtual Reality (IPI Press, USA) and the Editor in Chief of the Journal of WSCG. He has been a member of several international program committees of prestigious conferences and workshops. He is a member of ACM SIGGRAPH, IEEE and Eurographics Association. He became a Fellow of the Eurographics Association in 2010.
Prof. Vaclav Skala has published over 200 research papers in scientific journals and at international research conferences. His current research interests are computer graphics, visualization and mathematics, especially geometrical algebra, algorithms and data structures. Details can be found at http://www.VaclavSkala.eu

Prof. Rongjiang Pan is a professor in the School of Computer Science and Technology, Shandong University, China. He received a BSc in computer science, a MSc in computer science, a PhD in computer science from Shandong University, China in 1996, 2001 and 2005, respectively. During 2006 and 2007, he was a visiting scholar at the University of West Bohemia in Plzen under a program supported by the international exchange scholarship between China and Czech governments. He is now a visiting professor at the School of Engineering, Brown University from 2014 to 2015 under the support of China Scholarship Council.

He is a Member of the ACM. His research interests include 3D shape modeling and analysis, computer graphics and vision, image processing. He has published over 20 research papers in journals and at conferences.
Gravity control: modeling and experiments

Vitaly O. Groppen

Abstract—Paper describes statement and results of experiments which are based on the use of high voltage for verification of gravity control possibility. Used verification technology is based on the idea of substitution of energy distributed in the neighborhood above the upper surface of the plate-deployed capacitor by the material point with equivalent mass: force of the gravitational interaction of the plate with this point has opposite direction to the force of gravitational interaction of this plate with the Earth thus reducing this force. Results of experiments with different samples confirm the correctness of the proposed approach.

Keywords—Experimental verification, gravity control, high voltage, samples, weight measurement.

I. Introduction

The equation published in England by Sir Isaac Newton in 1667 [1] can be considered as the first attempt to describe the forces of gravity. About 250 years later, in 1915, in [2] Albert Einstein demonstrated a new theory of gravitation based on the Theory of Relativity. Movement of physical objects under the influence of high voltage discovered by Townsend Brown in 1921 cannot be considered as control of gravitational forces because this phenomenon is known to be caused by ionization of air atoms near acute and sharp edges [3]. The experiments outlined below are also based on the use of high voltage for verification of gravity control possibility. They are based on the idea of substitution of energy distributed in the neighborhood above the upper surface of the plate-deployed capacitor by the material point with equivalent mass: force of the gravitational interaction of the plate with this point has opposite direction to the force of gravitational interaction of this plate with the Earth therefore reducing the plate’s weight.

II. Main Principles

Electrodes of used in experiments plates are designed as metal strips of deployed capacitors so that their stored energy is distributed above the upper surface of a horizontally positioned plate (see Fig. 2a and Fig. 3a below). This energy $E_i$ of i-th plate is equal to:

$$\forall i: E_i = \frac{C_i U^2}{2},$$  (1)

where “$C_i$” is capacity of a charged plate-capacitor, “$U$” - power supply voltage.

The mass of this energy is determined as follows:

$$\forall i: m(E_i) = \frac{C_i U^2}{2c^2},$$  (2)

where “$c$” - velocity of light.

Below we use the model, in which:

a) the plates are disposed horizontally, so that the electrodes are on their upper surface;
b) distributed above the upper surface of the i-th plate energy is replaced by the equivalent body $D$, whose mass is determined by the expression (2).

Thus the force $F_i$ of the gravitational interaction between the i-th plate and body $D$ has a direction opposite to the force $F_e$ of gravitational interaction between this plate and the Earth (Fig. 1).

![Fig.1. The forces of interaction between the body D, i-th plate and Earth.](image)

Force $F_i$ value in accordance with (2), is determined by the low of gravity:

$$\forall i: F_i = \gamma \frac{m_i C_i U^2}{2R_i c^2},$$  (3)

where “$\gamma$” - gravitational constant, “$R_i$” - distance between the body D and the surface of the i-th plate (Fig. 1).

Denoting $F_e^0$ the weight of a plate before experiment whereas $F_e^1$ - its weight during experiment when the electrodes on its surface are applied to voltage equal to $U$, it is easy to determine value $F_i$: $\forall i: F_i = F_e^0 - F_e^1$.  (4)

Fixing, for each plate in the experiments according to (4) all components of the equation (3) except distance $R$, the latter for i-th plate can be defined as:
In a first approximation, the experimentally obtained values of $R_i(U)$ can be set by a polynomial:

$$\forall i: R_i(U) = \sum_{j=0}^{\gamma^2} k_{ij} U^j,$$  \hspace{1cm} (6)

where “$k_{ij}$” – $j$-th coefficient in equation (6) for $i$-th plate.

Substituting (6) into equation (3) it is possible to predict $F_i$ values outside the range of voltages $U$, fixed during the experiments, using the expression:

$$\forall i: F_i = \gamma \frac{m_i C U_j^2}{2 \sum_{j=0}^{\gamma^2} k_{ij} U^j} \frac{c^2}{2g},$$ \hspace{1cm} (7)

Using (7) it can be shown that the lifting force $F_i$, which is equal to the weight of the $i$-th plate, corresponds to the voltage $U$, being solution of the following square equation:

$$\forall i: \sum_{j=0}^{\gamma^2} k_{ij} U^j = \frac{U}{\gamma \cdot C} \frac{c^2}{2g}.$$  \hspace{1cm} (8)

III. Equipment, experiments statement and samples for gravity control

To verify the assumptions above during the first series of experiments (see [4], [5]) we used different light and thin fiber glass with lavsan cover plates with two groups of cooper electrodes on the upper side of each plate (Fig. 2a) connected to the high voltage power supply IVNR-20/10 (guarantying voltage range 1–20 kV, power 200 wt., see Fig. 2b,1) and installed on the weight table of precise electronic scale AV-60/01-S (precision is equal to 0.0001 g, maximum weight – 60 g., the settling time of weighting mode – about 10 minutes, Fig. 2b,2).

Fig. 2. Geometry of electrodes on the upper surface of the plates (a) and equipment (b) used in the first series of gravity control experiments.

Detailed descriptions of scheme and results of these experiments is provided in [5], there seems to be three typical sources of weight value mistakes during these series of experiments:

- experiments for direct weight measurement of samples under high voltage resulted in direct interaction of electronic circuit of the scale and its sensor with the electric field on the upper surface of a sample often resulting in distortions in indications of weight by the scale (Fig. 2b, 3) and even in blocking of electronics of the scale;
- aspiration to increase the energy stored by sample-capacitor increasing its capacity by reducing the distance between the electrodes, resulted in downturn of breakdown voltage, which ultimately contributed to the opposite effect - a decrease of stored energy;
- as shown in [5], any prolonged exposure of different samples based on fibre glass with lavsan cover to high voltage leads to its’ electrical breakdown, thus eliminating possibility of a repetition of the experiment.

To minimize the errors indicated above, during the second series of experiments we used:

- new samples with better resistance to electrical breakdown made of granite with two spaced apart parallel copper strips, attached to the top of each granite rectangle (see below Fig. 3a and Table 1);
- new precision balance AB-200 with maximum weight equal to 200 gram and precision equal to 0.001 g (Fig 3b), which is not exposed to electromagnetic radiation;
- new scheme of experiments making use of all the possibilities of the new equipment.

To improve the accuracy of weighing along with the standard metal weights were used narrow strips of paper. The scheme of experiments for each sample was determined by the following algorithm.

Fig. 3. The copper electrodes on the upper surface of the granite plates (a) and mechanical scale AB-200 (b) used in the second series of experiments.
Table 1. Parameters of the samples used during the second series of experiments.

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<tr>
<td>1</td>
<td>Upper surface area</td>
<td>a</td>
<td>0.0063</td>
</tr>
<tr>
<td>2</td>
<td>Total surface area</td>
<td>a</td>
<td>0.0158</td>
</tr>
<tr>
<td>3</td>
<td>Weight</td>
<td>a</td>
<td>211.0</td>
</tr>
<tr>
<td>4</td>
<td>Thickness</td>
<td>a</td>
<td>10.0</td>
</tr>
<tr>
<td>5</td>
<td>Distance between the electrodes</td>
<td>a</td>
<td>30.0</td>
</tr>
<tr>
<td>6</td>
<td>Width of the electrodes</td>
<td>a</td>
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<tr>
<td>7</td>
<td>Capacity</td>
<td>a</td>
<td>6.166</td>
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<tr>
<td>8</td>
<td>Material of the plate basis</td>
<td>a</td>
<td>Granite</td>
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Algorithm

Step 1. A sample is installed on the left weighing pan of precision mechanical scale AB-200 and balanced.

Step 2. A new, previously unused voltage is selected by the power supply IVNR-20/10 and applied to the sample.

Step 3. As during the previous step scale was unbalanced, we restore balance by removing part of the weights from the right weighing pan.

Step 4. The high voltage source is turned off.

Step 5. The weights tucked away on the third step of the current iteration are weighed by precision electronic scales AV-60/01-S. The total weight of these weights is shown by the scale display.

Step 6. We fix voltage, selected on the second step of the current iteration and corresponding change of weight of used sample.

Step 7. If experiments are carried out with all the planned values of high voltage, then go to step 8, otherwise mechanical scales AB-200 are balanced again and we go to the second step.

Step 8. The series of experiments with the sample selected at the first step is completed.

IV. Results of experiments

Voltage and corresponding change of weight for each sample are presented in the appendix whereas equation (6) coefficients for each sample described in the Table 3 are presented below (distance $R_i$, $i \in \{a, b, c\}$, is given in meters, voltage $U$ – in volts):

$$R_a(U) = -2.998478 \cdot 10^{-14} + 7.847125 \cdot 10^{-18} \cdot U - 3.021905 \cdot 10^{-22} \cdot U^2; \quad (9)$$

$$13 \text{kV} \leq U \leq 16 \text{kV}.$$  

$$R_b(U) = 1.034423 \cdot 10^{-14} - 4.466534 \cdot 10^{-19} \cdot U + 1.060829 \cdot 10^{-23} \cdot U^2; \quad (10)$$

$$9 \text{kV} \leq U \leq 20 \text{kV}.$$  

$$R_c(U) = 1.554385 \cdot 10^{-15} - 1.094039 \cdot 10^{-19} \cdot U + 3.561179 \cdot 10^{-23} \cdot U^2; \quad (11)$$

$$2.5 \text{kV} \leq U \leq 15 \text{kV}.$$  

In the equations (9) and (10), the relative error does not exceed five percent, with the voltage changes in the range of $9 - 20$ kV, whereas in equation (11) the relative error does not exceed 29 percent if the voltage changes in the range of $2.5 - 15$ kV.

The relative deviation of the calculated lifting force obtained through the equation (9) or (10) substitution into equation (7) from its experimental value does not exceed ten percent whereas the upper bound of the relative deviation of the calculated lifting force resulting equation (11) substitution into equation (7) is equal to the 46%.

Substitution of (9) into (8) for voltage $U$ value determination in the case when the lifting force $F_a$ value is equal to the weight of the sample “a” in Fig. 3a, results in the values equal to either 1990.744 V or to 49843.05 V. Similar substitution of (10) into equation (8) for voltage $U$ value determination for the case when the lifting force $F_b$ value is equal to the weight of the sample “b” in Fig. 3a, results in the values equal to either 13396.5 V or to 72788.28 V. The control series of experiments with the lower values of voltage applied to the both samples did not result in the corresponding lifting forces. The efforts to predict constant voltage $U$ value resulting in lifting force equal to the weight of sample “c” in Fig. 3a, by substitution of (11) into equation (8) failed: the resulting
V. Conclusions

The above presented approach allows us to make the following conclusions:

1. Results of experiments convince of the correctness of the proposed idea of gravity control.
2. As the lower values of voltage applied to the samples “a” and “b” during the experiments did not result in the corresponding lifting forces and the efforts to predict constant voltage $U$ value resulting in lifting force equal to the weight of sample “c” using dependence distance / voltage defined by equation (6) failed, further experiments will be focused on:
   • a wider range of voltages, overlapping defined above upper voltage values;
   • the search of dependence distance / voltage different from the model defined by equation (6).

Appendix

The tables below contain the results of experiments with the samples "a", "b" and "c" whose parameters are given above in Table 1. The first column of each table contains the number of experiment of the current series of experiments, the second column - the voltage applied to the sample, and the third - the proper lifting force.

Table 2. Voltage and corresponding change of lifting force $F_a$ for the sample “a” (Fig. 3a)

<table>
<thead>
<tr>
<th>#</th>
<th>$U$ (V)</th>
<th>$F_a$ (kg)</th>
</tr>
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<td>1</td>
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<td>0</td>
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<tr>
<td>2</td>
<td>13000</td>
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</tr>
<tr>
<td>3</td>
<td>14000</td>
<td>0.0000231</td>
</tr>
<tr>
<td>4</td>
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</tr>
<tr>
<td>5</td>
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</tr>
<tr>
<td>6</td>
<td>15500</td>
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</tr>
<tr>
<td>7</td>
<td>16000</td>
<td>0.0000393</td>
</tr>
</tbody>
</table>

Table 3. Voltage and corresponding change of lifting force $F_b$ for the sample “b” (Fig. 3a)

<table>
<thead>
<tr>
<th>#</th>
<th>$U$ (V)</th>
<th>$F_b$ (kg)</th>
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</thead>
<tbody>
<tr>
<td>1</td>
<td>0</td>
<td>0</td>
</tr>
<tr>
<td>2</td>
<td>9000</td>
<td>0.0000144</td>
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<tr>
<td>3</td>
<td>10000</td>
<td>0.0000194</td>
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<td>4</td>
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Table 4. Voltage and corresponding change of lifting force $F_c$ for the sample “c” (Fig. 3a)

<table>
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<tr>
<td>8</td>
<td>15000</td>
<td>0.000009</td>
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</table>

Acknowledgment

I’m pleased to thank Ann Minevitch for assistance in preparing of English version of this work.

References


Numerical investigation of diffusion turbulent propane/air flame

S. Morsli, M. El Ganaoui, A. Sabeur-Bendehina

Abstract—The turbulent diffusion flames can be found in a wide variety of thermal energy production systems. The burners operating with this type of flames are used for example in combustion chambers. The need to understand the structure of these flames has been our motivation to carry out a numerical study of the turbulent aero thermochemistry of C3H8 / Air diffusion flames in a burner using numerical simulations. Due to the high temperature and velocity gradients in the combustion chamber, the effects of equivalence ratio (φ) and oxygen percentage (γ) in the combustion air are investigated for different values of φ between 0.5 and 1.0 and γ between 10 and 30%. In each case, combustion is simulated for the fuel mass flow rate resulting in the same heat transfer rate (Q). Numerical calculations are performed for all cases using computational fluid dynamics code (Fluent CFD code). The results shown that the increase of equivalence ratio corresponds to a significantly decrease in the maximum reaction rates and the maximum temperature increase with the increases of oxygen percentage. Mixing hydrogen with propane causes considerable reduction in temperature levels and a consequent reduction of CO emissions.

Keywords—Combustion, Computational methods, diffusion, turbulent.

I. INTRODUCTION

Combustion is defined as the burning of a fuel and oxidant to produce heat and/or work. It is the major energy release mechanism in the Earth and key to humankind's existence. Combustion includes thermal, hydrodynamic, and chemical processes [1].

It starts with the mixing of fuel and oxidant, and sometimes in the presence of other species or catalysts. The fuel can be gaseous, liquid, or solid and the mixture may be ignited with a heat source.

When ignited, chemical reactions of fuel and oxidant take place and the heat release from the reaction creates self-sustained process: Turbulent combustion of hydrocarbon fuels and the incineration of various industrial by-products and wastes are an integral part of many segments of the chemical process and power industries.

Research into transport phenomena in energy systems and applications has substantially increased during the past a few decades due to its diversity in applications. This makes the special issue a most timely addition to existing literature.

The primary objectives in burner design are to increase combustion efficiency and to minimize the formation of environmentally hazardous emissions, such as CO, unburned hydrocarbons (HC) and NOx. Critical design factors that impact combustion include: the temperature and residence time in the combustion zone, the initial temperature of the combustion air, the amount of excess air and turbulence in the burner and the way in which the air and fuel streams are delivered and mixed [2].

This work considers the combustion of propane with air due to the high temperature and velocity gradients in combustion chamber using a single burner element. In order to investigate the effect of Therefore, CFD codes can serve as a powerful tool used to perform low cost parametric studies. The CFD codes solve the governing mass, momentum and energy equations in order to calculate the pressure, concentrations, velocities and temperatures fields.

II. MATHEMATICAL MODEL

The purpose of this simulation is to study the effect of oxygen percentage on the combustion; the combustion of fuel with air is examined at various oxygen percentages in the air by using Fluent CFD code in the combustion chamber of a burner. For this case, the combustion of propane (C3H8) with air, in a burner was considered. The two-dimensional axisymmetric model and geometric configuration of the burner are shown in Figure 1. As apparent from this figure, the fuel and air inlets are coaxial and merge downstream. It is assumed that the burner wall is under ambient conditions and that the walls near the air and fuel inlets are isolated. The model used for the numerical calculations include the RNG(renormalization group theory) k–ε for turbulent flow however for chemical species transport and reacting flow, the eddy-dissipation model with the diffusion energy source option is adopted. The mixture (propane-air) is assumed as an ideal gas; no-slip condition is assumed at the burner element walls.
The governing equations for mass, momentum and energy conservation, respectively, for the two-dimensional steady flow of an incompressible Newtonian fluid are:

\[ \frac{\partial \rho}{\partial t} + \rho \frac{\partial u_i}{\partial x_i} = 0 \]  \hspace{1cm} (1)

And the mass conservation equation for \( k \) species is written as follows:

\[ \frac{\partial \rho Y_k}{\partial t} + \frac{\partial (\rho (u_i + V_{k,i}) Y_k)}{\partial x_i} = \dot{\omega}_k \]  \hspace{1cm} (2)

For \( K = 1; N \) and \( V_{k,i} \), is the \( i \) component of the diffusion velocity \( V_k \) of the \( k \) species and \( \dot{\omega}_k \) is the production rate of the \( k \) species. We defined also the momentum conservation equation

\[ \frac{\partial \rho u_i}{\partial t} + \frac{\partial (\rho u_i u_j)}{\partial x_j} = -\frac{\partial P}{\partial x_i} + \frac{\partial \tau_{ij}}{\partial x_j} + \rho \sum_{i=1}^{N} f_{k,i,j} \]  \hspace{1cm} (3)

And

\[ \tau_{ij} = -\frac{2}{3} \mu \frac{\partial u_i}{\partial x_j} + \mu \left( \frac{\partial u_i}{\partial x_j} + \frac{\partial u_j}{\partial x_i} \right) \]  \hspace{1cm} (4)

The energy conservation equation is defined as:

\[ \frac{\partial \rho H}{\partial t} + \frac{\partial (\rho u_i h_i)}{\partial x_i} = \dot{\omega}_T + Q + \frac{\partial D P}{\partial t} + \frac{\partial}{\partial x_i} \left( \lambda \frac{\partial T}{\partial x_i} \right) + \frac{\partial}{\partial x_i} \left( \rho \sum_{k=1}^{N} V_{k,i} Y_k h_{k,i} \right) \]  \hspace{1cm} (5)

Two additional equations for the RNG \( k-\varepsilon \) turbulence model- The turbulence kinetic energy, \( k \), and the dissipation rate, \( \varepsilon \), are determined using the following transport equations, respectively:

\[ \frac{\partial \rho (k Y_k)}{\partial t} + \frac{\partial (\rho (k u_i Y_k))}{\partial x_i} = \frac{\rho}{\rho_x} \left( \frac{\partial \alpha_i}{\partial x_i} \frac{\partial k Y_k}{\partial x_i} \right) + G_k - \rho \varepsilon \]  \hspace{1cm} (6)

\[ \frac{\partial \rho (\varepsilon Y_k)}{\partial t} + \frac{\partial (\rho (\varepsilon u_i Y_k))}{\partial x_i} = \frac{\rho}{\rho_x} \left( \frac{\partial \alpha_i}{\partial x_i} \frac{\partial \varepsilon Y_k}{\partial x_i} \right) + (C_{1\varepsilon} G_k - C_{2\varepsilon} \varepsilon - \chi) \]  \hspace{1cm} (7)

This model, based on the work of Magnussen and Hjertager, called the eddy-dissipation model. The net rate of production of species \( i \) due to reaction \( r \), \( R_{i,r} \), is given by the smaller of the two expressions below:

\[ R_{i,r} = v_{i,r} W_{\text{air}} A \rho \frac{\varepsilon}{k} \min \left( \frac{Y_{E}}{v_{r,R} W_{\text{air}}}, \frac{\sum_{i} \rho Y_{P}}{\sum_{i} v_{i,r} W_{\text{air}}} \right) \]  \hspace{1cm} (7)

B. Boundary conditions

- At the fuel inlet \((x = 0 \text{ and } 0 < r < r_f)\), \( u_x = U_f \), \( U_r = 0 \), and \( T = T_{\text{in}} \)
- At the air inlet \((x = 0 \text{ and } r_i < r < r_0)\), \( u_x = U_{\text{air}} \), \( U_r = 0 \), \( T = T_{\text{in}} \)
- At the isolated walls \((x = 0, r_f < r < r_i \text{ and } r_0 < r < R)\), \( \frac{\partial T}{\partial x} = 0, U_r = 0 \)
- At the burner wall \((r = R, \text{ and } 0 < x \leq L)\)

III. COMBUSTION AND REACTION MECHANISM

The simplest description of combustion is that it is a process that converts the reactants available at the beginning of combustion into products at the end of the process. The most common combustion processes encountered in engineering are those which convert a hydrocarbon fuel (which might range from pure hydrogen to almost pure carbon (C), e.g. coal) into carbon dioxide (CO\(_2\)) and water (H\(_2\)O). In this study, the combustion of fuel with propane is modeled with a one-step reaction mechanism (NR=1). The reaction mechanism takes place according to the constraints of chemistry, and is defined by:

\[ C_5H_8 + \frac{2}{\phi} O_2 + \frac{5}{\phi} \frac{100 - \gamma}{\gamma} O_2 \rightarrow 3CO_2 + 4H_2O \]  \hspace{1cm} (8)

\[ \frac{5}{\phi} \frac{100 - \gamma}{\gamma} N_2 + \frac{5}{\phi} \frac{100 - \phi}{\phi} O_2 \]

Where: \( \phi = 5W_{O_2} + (100 - \gamma) / \gamma W_{N_2} m_{\text{fuel}} \)

\[ / (W_{\text{fuel}} m_{\text{air}}) \]  \hspace{1cm} (9)
Fluent 6.3.26 [3] was chosen as the CFD computer code which uses a finite-volume procedure to solve the Navier–Stokes equations of fluid flow in primitive variables such as velocity, and pressure. It can also model the mixing and transport of chemical species by solving conservation equations describing convection, diffusion, and reaction sources for each component species. The RNG model was used as a turbulence model in this study [4]. The constants are resumed in Table I. While the simulation value and inlet velocities of air are presented in Table II and III.

The solution method for this study is axisymmetric. In order to achieve higher-order accuracy at cell faces, second-order upwinding is selected.

<table>
<thead>
<tr>
<th>Table I: Constant of Turbulence Model</th>
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<td>$C_\mu$</td>
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<td>$C_{1\varepsilon}$</td>
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<td>$C_{2\varepsilon}$</td>
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<td>R</td>
</tr>
<tr>
<td>L</td>
</tr>
<tr>
<td>T</td>
</tr>
<tr>
<td>$h_{amb}$</td>
</tr>
<tr>
<td>$P_{sp}$</td>
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<tr>
<td>$\rho_{air}$</td>
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<p>| Table III: Inlet velocities of air for $U_f=2.247$m/s and $Q=1000$W |
|-----------------|-----------------|-----------------|</p>
<table>
<thead>
<tr>
<th>$U_{air}$ [m/s]</th>
<th>$\phi = 0.5$</th>
<th>$\phi = 0.7$</th>
<th>$\phi = 1$</th>
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<tr>
<td>10</td>
<td>56.436</td>
<td>40.312</td>
<td>20.218</td>
</tr>
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A. Grid independence

Grid tests were adopted to ensure grid independence of the calculated results. So, the total cell number of 15000 cells was adopted. The grid distributions are uniform within each region. (Figure 2)

B. Result and Discussion

Reactions Rates

Figure 3 show the contours of reaction rates in the combustion chamber for the cases of $\phi = 0.5$, 0.7, 1.0, and oxygen percentage in air $\gamma = 10$, 20 and 30% respectively. The effect of reactions is apparent from these figures. We can see that with the increase of $\gamma$ these regions contract in the axial direction however they expand in the radial direction. So, the reaction rate decreases significantly with the increase of $\phi$.

The Mass Fraction of Species

The mass fraction of species (O$_2$, C$_3$H$_8$ and H$_2$O) in the combustion chamber, which are related to the distributions of reaction rates, are plotted in figure 4 for $\phi=0.5$, 1 and $\gamma=10\%$. The results obtained from this figure show that in the cases of $\phi < 1$, complete combustion occurs, while in the case of $\phi = 1$ it is very close to the complete combustion state.
For ($\phi = 0.5$ and $\gamma = 10\%$)

For ($\phi = 0.5$ and $\gamma = 30\%$)

For ($\phi = 0.7$ and $\gamma = 20\%$)

For ($\phi = 1.0$ and $\gamma = 10\%$)

For ($\phi = 1.0$ and $\gamma = 20\%$)

For ($\phi = 1.0$ and $\gamma = 30\%$)

Fig. 3 Contours of reaction rates for different values of equivalence ratio $\phi$ and oxygen percentage in air $\gamma$.

For ($\phi = 1.0$ and $\gamma = 10\%$)

For ($\phi = 1.0$ and $\gamma = 20\%$)

For ($\phi = 1.0$ and $\gamma = 30\%$)

Fig. 4 Mass fractions of species $\phi = 0.5, 1$ and $\gamma = 10\%$
Temperature distribution

Knowing that the heat, which is released in the chemical reaction and transferred into the flowing gas including the reactant and product molecules [5], increases the temperature of this gas. In order to view the temperature distribution, the variations of temperature along the axis of burner are plotted in figure 5.

It is very important to notice that the increase of $\phi$ significantly reduces the gradients so, larger temperature gradients occur in the axial direction (especially between $x = 0$ and about 0.20 m). The heat calculations performed for the each fuel case bring out that in the case of $\phi < 1$, the total heat per unit mass released in the combustion,

Consequently, the results obtained from this figures show that in the cases of $\phi < 1$, complete combustion occurs, while in the case of $\phi = 1$ it is very close to the complete combustion state.

VI. CONCLUSION

The combustion of propane with air in a burner was considered and the effect of the equivalence ratio and oxygen percentage in air investigated, for different numerical values. The specific conclusions derived from this study can be listed briefly as follows:

- The increase of $\phi$ reduces significantly the reaction rate levels.
- In the case of $\phi < 1$, the complete combustion occurs, and the combustion in the case of $\phi = 1$ is very close to the complete combustion state.
- The maximum temperatures in the combustion chamber increase with the increases of $\gamma$ (from 10 to 30%) and $\phi$ (from 0.5 to 1.0), respectively.

Consequently, the results of this study clearly demonstrate that mixing hydrogen with propane causes considerable reduction in temperature levels and a consequent, reduction of CO.

REFERENCES

NOMENCLATURE

Latin symbols

<table>
<thead>
<tr>
<th>Symbol</th>
<th>Definition</th>
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<tbody>
<tr>
<td>CFD</td>
<td>Computational fluid dynamics</td>
</tr>
<tr>
<td>Cμ, Cε1, Cε2</td>
<td>Coefficients in k–ε turbulence model</td>
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<tr>
<td>$f_{k,i}$</td>
<td>species $K$ in the direction $i$</td>
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<td>$\mu$</td>
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<td>$h_{r,k}$</td>
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<td>Pressure</td>
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<td>$Q$</td>
<td>Heat transfer rate</td>
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<td>$\eta_o$</td>
<td>Outer radius of air inlet</td>
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<tr>
<td>$T$</td>
<td>Temperature</td>
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<td>$Y_k$</td>
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Greek symbols

<table>
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<tr>
<td>$\delta_{ij}$</td>
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</tr>
<tr>
<td>$\phi$</td>
<td>Equivalence ratio</td>
</tr>
<tr>
<td>$\lambda$</td>
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<td>Oxygen percentage in air</td>
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<tr>
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Abstract—Up to now, the Chilean Energy system has fulfilled the energy needs of Santiago de Chile considerably well. However, development trends of the current system impose significant future risks on the energy system. A detailed sustainability analysis of the energy sector of the Metropolitan Region of Santiago de Chile was conducted, using selected energy indicators and a distance-to-target approach. Risks for the sustainable development of the energy sector are detected, such as increasing concentration in the energy sector, import dependency for fossil fuels and increasing CO₂ emissions from energy production. Options towards a more sustainable development of the Megacity of Santiago within the national Chilean energy system are assessed, such as the enhancement of energy efficiency and an increased use of renewable energies.

Keywords—future energy situation, indicators of sustainability, scenarios, sustainable development in developing countries

I. INTRODUCTION

Chile’s energy system is characterized by a strong growth as well as a high degree of privatization and economic concentration of energy service providers, which are controlled by the National Energy Commission (CNE) and the Energy Ministry. The total energy use in Chile (primary consumption) increased from 789 PJ (Petajoule) in 1999 to 1,593 PJ in 2012 [1], an annual average growth of 5.6%. This surge in energy consumption is for the most part a result of population growth, a highly dynamic economic development, and deficiencies in the effective use of energy resources. Approximately three quarters of the primary energy is based on fossil energy resources, which are almost completely imported from abroad.

The energy consumption of the Santiago Metropolitan Region (SMR) is dominated by the traffic sector accounting for 38% of final energy. This is followed by industry (27%), households (22%), and trade and services (13%) (data based in own calculation on [2]). By taking a closer look at the distribution of the final use of energy in terms of energy sources, it can be seen that – especially due to the energy use in the traffic sector – oil derivatives account for almost half of the final energy consumption. The other half is more or less equally distributed between gas and electricity.

The SMR can cover only 25% of its electricity demand on its own; therefore 3/4 of the electricity has to be imported. In the SMR 50% of the power is generated by fossil-fired thermal power plants, 50% by hydroelectricity.

Industry consumes the largest share of the electricity in the SMR (30%), followed by households (26%), and trade (22%). Mining accounts for 7%, agriculture for 2%, and the other sectors for 13% [3]. The power consumption per capita in the category “households” varies considerably in the different districts of Santiago. While people in Vitacura – one of the municipalities with the highest average household income – use around 1,200 kWh electricity per capita and year, people in the “poorer” municipalities Alhué and El Monte only consume less than 360 kWh [4].

II. ENERGY SCENARIOS

In a first step two framework scenarios for the years 2030 have been developed and the interrelations between the framework scenarios and the energy sector were analyzed. This was the starting point for the development of energy scenarios (according to [5]).

The future development of the indicators was estimated and assessed in the framework of these energy scenarios, partly based on the MESAP/PlaNet model (see [6]). The parameters which were included in the modeling are presented in table 1.

<table>
<thead>
<tr>
<th>Area</th>
<th>BAU (Business as usual)</th>
<th>CR (Collective Responsibility)</th>
<th>MI (Market Integration)</th>
</tr>
</thead>
<tbody>
<tr>
<td>Role of hydropower</td>
<td>Realization of the large-scale plant HidroAysen</td>
<td>Focus on small hydroelectric power plants</td>
<td>Realization of the large-scale plant HidroAysen</td>
</tr>
<tr>
<td>Role of non-conventional renewable energy carriers</td>
<td>Implementation of agreed target values (5% for 2010 / 10% for 2024)</td>
<td>Strong increase of combined heat and power, wind, solar, geothermal energy, biomass</td>
<td>Intensified increase in power generation by wind and solar panels</td>
</tr>
<tr>
<td>Role of fossil fuels</td>
<td>Further investments in fossil fuel power plants</td>
<td>Gas as backup for renewable energies</td>
<td>Further investments in fossil fuel power plants</td>
</tr>
<tr>
<td>Transport sector: share of electric vehicles in all cars registered</td>
<td>6%</td>
<td>10%</td>
<td>10%</td>
</tr>
</tbody>
</table>

III. STAKEHOLDER INVOLVEMENT

In the course of the framing project “Risk Habitat Megacity”, a joint German-Chilean research initiative [8], 30 experts from research, federal and regional authorities, and the industry were interviewed. These interviews provided important information on the data basis and the current situation, but also on the assessments of the future of the energy supply in Santiago and Chile and thus also for
the further development of the scenarios.

A stakeholder workshop was held at the Economic Commission for Latin America and the Caribbean (ECLAC) in Santiago. The aim of the workshop was to discuss indicators and target values as well as scenarios. During the workshop, the 20 participants from universities, authorities, and the industry entered into – in some cases very lively – discussions, especially concerning the indicators.

One of the results of the participation process was, that nuclear power is no option for Chile because

a. the high risk of disastrous earthquakes,

b. the huge investment and follow up costs of the technology and

c. the very high national potential on renewable energy with very low follow up costs.

The energy results were presented in a concluding workshop. Here especially the presented options for action were topics for discussion.

In addition, the results were included into the dialog of the regional government of the metropolitan region (GORE) for the development of a regional “política de energías limpias”.

IV. SUSTAINABILITY ANALYSIS IN THE ENERGY SECTOR

A comprehensive set of 44 indicators was compiled on this basis of an analyses of international literature (see [9]). Then 16 “core sustainability indicators” were chosen from this list, especially regarding the appropriate representation of the rules of sustainability of the integrative concept of sustainable development [10]. Finally the following 8 indicators for the assessment of scenario results were chosen primarily based on the criteria (a) possibility to determine target values for the indicators and (b) availability of SMR data:

1. Share of rural households with no access to electricity
2. Duration of electricity supply interruptions
3. Total primary energy consumption
4. Energy intensity estimated as: Energy per GDP
5. Share of renewable energies in electricity production
6. Energy-related CO₂ emissions per capita
7. Energy import dependency estimated as: Percentage of primary energy use based on imported energy
8. Degree of economic concentration in the energy sector

Where no sufficient data was available for the SMR, national or regional values were collected for the respective indicators. Target values were identified for all indicators based on existing local, regional, or national values; if this was not possible, scientific expertise was taken as a basis.

Table 2 Scenario results ☺ = target reached, ☹ = more to do

<table>
<thead>
<tr>
<th>Indicator</th>
<th>Unit</th>
<th>Current value</th>
<th>Target value 2030</th>
<th>BAU 2030</th>
<th>CR 2030</th>
<th>MI 2030</th>
</tr>
</thead>
<tbody>
<tr>
<td>Share of rural households with no access to electricity in the SMR</td>
<td>%</td>
<td>0.4</td>
<td>0.0</td>
<td>☺</td>
<td>☹</td>
<td>☹</td>
</tr>
<tr>
<td>Duration of electricity supply interruptions in SMR</td>
<td>h</td>
<td>2.8</td>
<td>0.0</td>
<td>☹</td>
<td>☺</td>
<td>☺ 0.4</td>
</tr>
<tr>
<td>Total primary energy consumption</td>
<td>PJ</td>
<td>497 /2007</td>
<td>As low as posible</td>
<td>☹ 880</td>
<td>680</td>
<td>824</td>
</tr>
<tr>
<td>Energy per GDP</td>
<td></td>
<td>85 /2008</td>
<td>≤ 40</td>
<td>☹ 65</td>
<td>50</td>
<td>60</td>
</tr>
<tr>
<td>Share of renewable energies in electricity production</td>
<td>%</td>
<td>55 /2007</td>
<td>≥ 70</td>
<td>☺ 66</td>
<td>75</td>
<td>76</td>
</tr>
<tr>
<td>Energy-related CO₂ emissions per capita</td>
<td>t</td>
<td>4.8 /2007</td>
<td>2.5</td>
<td>☺ 5.3</td>
<td>3.1</td>
<td>4.6</td>
</tr>
<tr>
<td>Percentage of primary energy use based on imported energy</td>
<td>%</td>
<td>77 /2007</td>
<td>≤ 50</td>
<td>☹ 69</td>
<td>52</td>
<td>67</td>
</tr>
<tr>
<td>Degree of economic concentration in the energy sector</td>
<td>%</td>
<td>90 /2003</td>
<td>≤ 70</td>
<td>☹ 90</td>
<td>86</td>
<td>95</td>
</tr>
</tbody>
</table>
V. DISCUSSION OF RESULTS ACROSS SCENARIOS

The sustainability analyses of the current status and the scenario development of the indicators show a heterogeneous picture for the different scenarios. There are numerous positive developments, but also cases with large discrepancies from the target values have to be stated.

As an example, we will have a closer look at the indicator “percentage of primary energy use based on imported energy”.

One detected sustainability deficit for Chile (and the SMR) is the high dependency on fossil fuel imports, as this increases vulnerability to supply shortages and price increases. For many years Chile has been importing more than two thirds of the required energy resources. The import quota of primary energy increased from 50% in 1990 to more than 70% today. Reasons for this are the growing energy demand and the scarce reserves of conventional energy carriers in Chile.

Especially the import of natural gas from Argentina increased considerably in the past. In 2006/2007 Argentina dramatically reduced the export of natural gas to Chile by a breach of contract; this caused serious problems for Chile’s industry and private households. As a consequence, gas was to a large extent replaced by diesel and coal, which lead to a triplication of diesel imports. Also the import of oil and coal increased significantly during this period. Until 2010, approximately 1 billion US dollars were invested into two large international ports (Quintero and Mejillones) with a capacity of ca. 16 million m³ liquid gas (LNG) per day to compensate for the gas supply shortage from Argentina.

The scenario results show that different developments are possible here. Indeed the percentage share of the dependency on imports is being reduced in the BAU and MI scenarios. But the growing energy demand raises the absolute amount of imported energy raw materials. An absolute as well as a relative decrease would only be possible in CR resulting from both the significantly smaller increase of the energy demand and the stronger use of domestic renewable energies (see fig. 1) without realizing the large-scale plant HydroAisen. In this scenario a reduction of the dependency on imports to 55% could be nearly achieved.

VI. ENERGY POLICY FOR MORE SUSTAINABILITY

During the last years, Chile took some energy-political steps towards a more sustainable energy supply. An Energy Ministry was created on the national level, a program for energy efficiency was launched, and the regional government responsible for the SMR is working on a “política de energía limpia” for the metropolitan region. In a future according to the CR scenario it would be possible to achieve some or come near to the sustainability target values for the selected indicators, e.g. for the indicator dependency on imports. In a world according to the BAU or MI scenario this would be different. From the foresight aspect it is therefore necessary to take further action.

One important measure is to enhance energy efficiency beyond the current activities in Chile to achieve the objective of an efficiency increase of 20% by 2020. In addition, the course of the development of thermal energy production plants, gas-fired power plants should be preferred to coal-fired power plants since these can be controlled more flexibly and are therefore better suited for a continuous energy supply in combination with power plants which are based on fluctuating renewable energy carriers dependent on climatic conditions.

All in all, numerous studies conclude that non-conventional renewable energies (NCRE) in Chile have enormous potential for the future of the energy sector [12], [13] by far exceeding current policy targets. Universitas of Chile and Technical University Federico Santa Maria estimate, that by 2025, non-conventional renewables could account for up to 30 TWh produced by almost a 6 GW installed capacity and provide more than 30% of the total power generation [14].

The enormous potential for the use of energy from renewable sources is not equally distributed over Chile. The main reserves of energy potentials can be found for

- hydropower in the Andes
- biomass in the forests in the south, in the agricultural areas in the Central Region and the waste of the metropolitan region,
- wind power in the south,
- solar energy in the north,
- ocean energy in the west and
- geothermal energy all over the country.

To ensure that the energy from the different regions arrives in the SMR, the infrastructure of supply lines has to be improved. Concerning the large potential of solar energy available in the north, there is the possibility of producing hydrogen or methane from solar electricity and feeding these gases into the natural gas grid which supplies the different regions of Chile coming from Argentina but is not used to capacity due to the Argentine supply shortage.

The search for and analysis of potentials from renewable energies should be started in the SMR since transport-related line losses are increasing with growing distance to the energy source. An evaluation of the suitability of existing roof areas for generating solar energy based on data...
from an “overflight” of the SMR would be quite easy to realize, as it is currently done in Osnabrueck (Germany) [15], but also waste material and sludge from the SMR have high potentials for a regional supply with renewable energy.

To achieve the necessary goals of increasing both the energy efficiency and the share of renewable energies in the total energy production, not only technological measures are required but also the establishment of political framework conditions which hardly exist in Chile and the SMR today. Here especially the introduction of relevant taxes and charges for the use of fossil energy carriers has to be mentioned as well as the fact that the feed-in of renewable energies into the grid should be facilitated and given priority.

Another important measure will be the analysis of the local energy resources (hydro, biomass, solar, ocean, and geothermal) of the SMR area and the two neighboring regions V and VI and to public these information for free in the internet. This could be the basis for investors to plan their energy projects for the SMR.

REFERENCES

A scalable method for efficient stem cell donor HLA genotype match determination

D. Georgiev, L. Houdová, M. Fetter, and P. Jindra

Abstract—Finding suitable stem cell donors comprises three independent processes: donor pool HLA typing, donor HLA haplotype inference, and search for donor HLA genotype matches. For practical and technical reasons, these processes are often decoupled leading to informational losses along the way. A method is presented that eliminates some of the technical challenges by considering all three steps together. The method relies on two practical assumptions: there exists a set of common haplotypes and the matched target is typed at high resolution. Under these two assumptions, sufficient statistical HLA analysis of a stem cell donor pool is performed to identify donors that most likely match a given genotype. The presented common haplotype Expected-Maximization (chEM) method is scalable in the number of loci, the number of alleles, and typing ambiguity, overcoming the known curse of dimensionality for the problem of HLA haplotype inference. The practical value is demonstrated on real world data provided by the Czech National Marrow Donor Registry. It is shown the chEM method significantly reduces the field of potential matches when compared to an existing match algorithm.

Index Terms—stem cell donors, HLA refinement, statistical methods

I. BACKGROUND

A brief review of the typing, resolution, and matching processes is given. For the sake of simple exposition, unnecessary details are omitted. For more information see [1].

A. Typing methods and nomenclature

Donors are usually typed for up to five loci along chromosome 6: HLA-A/B/C/DRB1/DBQ1. Typing results of donors in a given registry bare different levels of resolution, based on time and location of recruitment. Low level methods rely on antibody-based serological tests. Intermediate resolution molecular methods are based on hybridisation with sequence-specific oligonucleotide probes or PCR amplification with allele specific primers. High resolution methods are based on DNA sequencing tools, which aim to return a four digit number, e.g., A*03:01 or A*03:26, describing the exact amino acid sequence of the corresponding allelic coding region. Ambiguity regarding the amino acid sequence is possible for highest resolution typed donors if their chromosome pair contains an indeterminate allele combination. Decreasing the typing resolution further increases possible allele ambiguity.

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The above donor is heterozygous at the four typed loci, where the loci HLA-A/B/DRB1 are typed at intermediate resolution and HLA-C is typed at low resolution. The locus HLA-DQB1 is not typed implying it has 509 possible alleles on each chromosome. Such a donor has (at high resolution) \( 1.9 \times 10^{10} \) possible haplotypes, \( 5.3 \times 10^{17} \) possible HLA genotypes, and \( 4.3 \times 10^{18} \) possible haplotype pairs. If the typing data is truncated to low resolution, then the ambiguity is reduced significantly: 224 possible haplotypes, 196 possible genotypes, and 1568 possible haplotype pairs. If, in addition, the haplotype is limited to the typed loci, then (at low resolution) there are 16 possible haplotypes, 1 possible genotype, and 8 possible haplotype pairs.

B. Statistical methods for HLA refinement

Additional typization to increase HLA resolution and reduce the ambiguity (at high resolution) is rarely implemented for a large pool of donors. Instead statistical methods are used to refine a given HLA typing. In the simplest approach, allele frequencies are calculated independently at each locus and used to compute genotype probabilities without accounting for linkage effects. Real world populations, however, share common ancestry and exhibit nonuniform mating patterns that lead to linkage disequilibrium manifested in allele correlations across HLA haplotypes [3].

There is no gold standard for probabilistic haplotype modelling generally accepted by the stem cell donor registries. Three types of methods are available [1]. The Clark’s algorithm approaches haplotype modelling parsimoniously beginning with the list of homozygous haplotypes and then

<table>
<thead>
<tr>
<th>A</th>
<th>B</th>
<th>C</th>
<th>DRB1</th>
</tr>
</thead>
<tbody>
<tr>
<td>68:</td>
<td>27:</td>
<td>07(348)</td>
<td>11:</td>
</tr>
<tr>
<td>JAH(16)</td>
<td>HMHK(13)</td>
<td>HMXW(12)</td>
<td></td>
</tr>
<tr>
<td>03:</td>
<td>07:</td>
<td>01(90)</td>
<td>15:</td>
</tr>
<tr>
<td>JAGH(50)</td>
<td>HMGA(31)</td>
<td>HMYC(17)</td>
<td></td>
</tr>
</tbody>
</table>

The above donor is heterozygous at the four typed loci, where the loci HLA-A/B/DRB1 are typed at intermediate resolution and HLA-C is typed at low resolution. The locus HLA-DQB1 is not typed implying it has 509 possible alleles on each chromosome. Such a donor has (at high resolution) \( 1.9 \times 10^{10} \) possible haplotypes, \( 5.3 \times 10^{17} \) possible HLA genotypes, and \( 4.3 \times 10^{18} \) possible haplotype pairs. If the typing data is truncated to low resolution, then the ambiguity is reduced significantly: 224 possible haplotypes, 196 possible genotypes, and 1568 possible haplotype pairs. If, in addition, the haplotype is limited to the typed loci, then (at low resolution) there are 16 possible haplotypes, 1 possible genotype, and 8 possible haplotype pairs.
expanding the list in an ad hoc manner until all donor genotypes are resolved. Expected-Maximization (EM) methods consider all possible genotype deconstructions and look for haplotype probabilities of the general population that maximise the likelihood of drawing the donor pool. Bayesian methods rely on more detailed models of population genetics to stochastically generate population haplotypes and empirically estimate haplotype frequencies [4].

Choice of method is determined by several factors. Usually EM and Bayesian methods outperform the Clark’s algorithm. Furthermore, EM methods have better convergence properties, include simpler parametrizations, but scale poorly with the number of loci, the number of alleles, and the allele ambiguity. Bayesian methods potentially have lower computational complexity but lack the implementation simplicity of the EM methods. For this reason, most stem cell registries implement their own version of the EM method [5], [6].

Below is a summary of the standard EM algorithm used for haplotype modelling.

The haplotype probability model is iteratively computed as follows.

1) Haplotype set $H = \{h_1, ..., h_n\}$: let $d$ be the number of donors and construct a set of haplotypes by deconstructing the donor genotypes.

2) Initialization of $p_0(k) \ L_0, s, S, \epsilon$: initialise the probability $p_0(k)$ of haplotype $h_k$ appearing in the population and the likelihood function $L_0 = 0$. Set the iteration count $s = 0$, the maximum iteration count $S$, and the convergence threshold $\epsilon$.

3) Start: while $s \leq S$, repeat the following steps, otherwise, terminate the algorithm without convergence.

i) For $k = 1, ..., n$, count the occurrence $n_k$ of haplotype $h_k$.

$$n_k = \sum_{i=1}^{d} \sum_{j=1}^{n} (1 + \delta_{k_j}) P(h_k, h_j|g_i), \quad (1)$$

where $\delta_{k_j}$ is the delta function, $g_i = \{h_{i11}, h_{i12}\}, \{h_{i21}, h_{i22}\}, ..., \{h_{i11}, h_{i12}\}$ is the $i$th donor’s genotype deconstructed into all possible haplotype pairs, and

$$P(h_k, h_j|g_i) = \frac{p_s(k)p_s(j)}{P(g_i)}, \quad (2)$$

$$P(g_i) = \sum_{j=1}^{r} p_s(h_{ij1})p_s(h_{ij2}). \quad (3)$$

ii) For $k = 1, ..., n$, compute the next iteration of the haplotype probabilities and the likelihood function:

$$p_{s+1}(k) = \frac{n_k}{n}, \quad (4)$$

$$L_{s+1} = \sum_{i=1}^{d} \log (P(g_i)). \quad (5)$$

iii) Check the convergence criteria: if $|L_{s+1} - L_s| > \epsilon$, increment $s$ and return to (i), otherwise, terminate with convergence.

The limitation of the standard EM algorithm is in Step 1. Only a small percentage of donors are typed at high resolution. Deconstructing genotypes of donors typed at intermediate or low resolution into all possible haplotypes generates a prohibitively large haplotype set. Hence, computational complexity attributed to typing ambiguity for actual donor pools is overly prohibitive for simple deployment of the standard EM algorithm [1]. In practice, EM algorithms are either executed at low resolution, where there is little ambiguity (see example above), or executed heuristically to explore specific linkages in partial haplotypes.

C. Donor matching

A donor is a match for a given patient if they share the same genotype. With some exceptions, a patient is typed at high resolution and hence ambiguity arises almost exclusively on the donor side. Two steps are commonly used to resolve this ambiguity. A simple solution is to convert typing data to a lower resolution, e.g., antigen split groups, where there is little ambiguity and common population haplotypes may be used to predict the missing information. Alternatively, the donor genotype is matched only across the typed loci, generating so called 6/6, 8/8, and 10/10 matches. Such matches are of the boolean type, a donor is classified as a potential match if their partially genotype possibly equals the partial genotype of the patient. Effectively, use of probabilistic haplotype models for HLA refinement is ignored.

II. RESULTS

The results in this paper comprise a method that overcomes the greatest shortcoming of the EM algorithm, its computational complexity caused by typing ambiguity at high resolution. Below the common haplotype EM (chEM) method is introduced and its scalability demonstrated in a deployment on a large portion of CNMRD donors typed at various loci and various resolution levels. A five-loci probabilistic model of the Czech population is derived. In addition, the model and the introduced tools are shown to outperform existing intermediate resolution methods in donor matching.

Assumption. The chEM method is based on the following assumptions.

A1) The list of common haplotypes is known.
A2) Patients are typed at high resolution.

Assumption 1 is likely considering most newly discovered alleles are rare [7]. Sequencing methods have also become affordable enough to where Assumption 2 is now generally true.

A. Common haplotype probability model

The derived method differs from the EM algorithm approach described in Section I-B in the following fundamental way. The haplotypes are decomposed into those that are common, denoted simply by $H = \{h_1, ..., h_n\}$, and those that are rare, denoted by $H^R = \{h_1^R, h_2^R, \ldots\}$. Rare haplotypes
either contain rare alleles (possibly still unknown) or a rare combination of alleles and as a result have a much lower probability of arising in the population.

The EM algorithm is modified and implemented in the following way to efficiently compute the common haplotype probability model.

1) Haplotype set $H = \{h_1, ..., h_n\}$: the set of common haplotypes is given.
2) Initialization of $p_0(k), L_0, s, S, \epsilon, M$: Initialise $p_0$ and $L_0$ as above. Set the values of $s, S,$ and $\epsilon$ as above. Set $M$ to be the number of maximum allowable haplotypes per donor.
3) Construct the relevant donor pool $D = \{1, ..., d\}$: a donor in the general pool is deemed relevant as follows.
   i) Set the iteration count $a = 1$ and initialise the donor’s list of potential haplotypes $H_p$ to $H$.
   ii) Consider all alleles possible for the donor at the $a$th locus and exclude from $H_p$ any haplotype that does not share one of these alleles.
   iii) If the size of $H_p \leq M$, the donor is deemed relevant, otherwise, if $a$ is less than the number of loci, increment $a$ and return to the previous step. If $a$ is equal to the number of loci, the donor is deemed irrelevant.
4) Construct the set of common haplotype pairs $g_i^f, i \in \{1, ..., d\}$: for donor $i$, first select one of the pre-computed common haplotypes, then search for its pair. The search is done in the same way as when identifying relevant donors (Step 3), the exception being the possible alleles are now limited to the complementary chromosome copies.
5) Continue as in Step 3 of the traditional EM algorithm described above with $g^f$ substituted for $g$.

Limiting the EM algorithm to the common haplotype list reduces the computational complexity in a number of ways. First, model complexity, i.e., the number of potential haplotypes, is explicitly bounded. As a result, the computational complexity attributed to typing ambiguity (represented by the maximal number of donor haplotypes $M$) can be independently controlled. In the standard EM algorithm, the number of donor haplotypes grows with $M$ and the number of donors $d$. Hence to reduce the model complexity, one must consider either high resolution donors or limit the number of donors. Lastly, a relevant donor must not necessarily be typed at high resolution. Their genotypes must merely deconstruct into an acceptable number of common haplotypes. In the case study below, for $M = 150$ the number of relevant donors $d$ from CNMRD is equal to $23,819$ (54% of the entire database). In the standard EM algorithm, any donor typed at intermediate resolution or below for 2 out of the 5 loci would be deemed irrelevant. In the present, there are only $3,339$ (7.6% of the entire database) such donors in the CNMRD. Hence, limiting the model to the common haplotypes greatly increases the statistical relevance of the results for fixed donor ambiguity.

B. Matching genotypes

The chEM algorithm presented in Section II-A only computes the probability model for the common haplotypes. What isn’t yet clear is how this model is useful in matching donors to patients, especially when the patient has a rare haplotype. We start by again considering the donor and patient genotypes deconstructed into sets of all possible haplotype pairs $g_D$ and $g_P$, respectively.

The standard search for a donor with a genotype that likely matches the patient’s genotype involves processing the general donor pool and for each donor evaluating the conditional probability

$$P(g_P|g_D) = \frac{P(g_P \cap g_D)}{P(g_D)},$$

where the Hardy–Weinberg equilibrium is used to compute the probabilities on the right hand side of the equation. For this, however, we need a full probability model of all possible haplotypes in the donor pool. The chEM algorithm of Section II-A only computes the probabilities $P(h_k|\text{donor pool contains 0 rare haplotypes})$, written here simply as $P(h_k|0)$.

If the occurrences of different rare haplotypes in the donor pool are equally likely and statistically independent and $\xi$ is the probability the donor pool has at least one rare haplotype, then the true probability is given by

$$P(h_k) = P(h_k|0) (1 - \xi) + \sum_{i=1}^{\infty} P(h_k|i) \left( \frac{\xi}{1 + \xi} \right)^i.$$  

Under assumptions A1 and A2, the probability $\xi$ is small. Hence we can approximate the match probability in Equation 6 by its Taylor series approximation.

$$P(g_P|g_D) = \sum_{i=0}^{\infty} \frac{\xi^i}{i!} \frac{\partial^i P(g_P|g_D)}{\partial \xi^i} \bigg|_{\xi=0}$$

Three scenarios are possible and characterised by the dominant terms in the above Taylor series.

Scenario 1) If $g_P$ includes at least one common haplotype pair, then

$$P(g_P|g_D) \approx \sum_{(h_1, h_2) \in g_P \cap g_D} P(h_1|0)P(h_2|0).$$

Scenario 2) If $g_P$ includes no common haplotype pairs and at least one pair with one common haplotype, then

$$P(g_P|g_D) \approx \sum_{(h_1, h_2) \in g_P \cap g_D} P(h_1|0)P(h_2|1) + P(h_1|1)P(h_2|0).$$

Scenario 3) If $g_P$ includes only rare haplotype pairs, then

$$P(g_P|g_D) \approx \sum_{(h_1, h_2) \in g_P \cap g_D} P(h_1|1)P(h_2|1).$$

The above facts suggest that to approximate the donor match probabilities, all rare haplotypes would have to be appended one at a time to the haplotype model. This is true if the probabilities themselves are of interest, however, for the purposes of finding the best donor, only the ordering of the probabilities is required.
Note, for any patient typed at high resolution, it is true that either a donor’s possible haplotype pairs include ALL of the patient’s possible haplotype pairs or NONE of the patient’s haplotype pairs. Correspondingly divide the general donor pool into the potential matches denoted by $D_P$ and the rest, where any donor $i \in D_P$ satisfies $g_i \supseteq g_P$ and any donor $j \notin D_P$ satisfies $g_j \cap g_P = \emptyset$.

Any two donors $i, j \in D_P$ satisfy the equality $g_i \cap g_P = g_j \cap g_P$. As a result, the numerator of the match probability approximations in the three scenarios is constant for all donors in $D_P$. This leads to the following fact regarding ordering of matching donors.

**Fact 1.** Consider the set $D_P$ of donors possibly matching a given patient’s genotype, represented by the set of haplotype pairs $g_P$. As the probability $\xi$ of a rare haplotype in the donor pool approaches 0, for any two donors $i, j \in D_P$,

$$P(g_P|g_i) > P(g_P|g_j) \text{ if and only if } P(g_i|0) < P(g_j|0),$$

where $P(g_i|0)$ and $P(g_j|0)$ are the genotype probabilities conditioned on there being no rare haplotypes in the donor pool.

The above fact states that only the common haplotype model is required to order the potentially matching donors. The match probability can be estimated using the common haplotype model only when the patient genotype can be reconstructed using common haplotypes. Otherwise, the haplotype model would have to be expanded to include all rare haplotypes.

### C. CNMRD derived haplotype model of the Czech population

The haplotype model was derived for a CNMRD donor pool satisfying the ambiguity criterion introduced in the chEM algorithm. CNMRD includes 44,256 donors, of which 39,403 (90%) are typed for HLA-A/B/DRB1 loci and 70% have some serological typization. Donor ambiguity was limited to at most 150 possible common haplotypes, met by 54% of the registry’s donors. Donors deemed overly ambiguous to yield significant statistical information included many that were typed for only two loci, e.g., HLA-A/B, and had up to 20 thousand possible common haplotypes.

The list of common haplotypes making up the derived haplotype model included 64,800 entries taken from [8]. This was assumed to be the list of all common haplotypes called for by Assumption 1. In deployment, however, the possible genotypes of numerous donors from the selected donor pool could not be reconstructed using the listed common haplotypes. All such donors had serologically typed HLA-C. Consultation with field experts revealed that older serological typing methods for HLA-C have a 30% error rate. Once HLA-C typing was ignored for these donors, common haplotype pairs were successfully used to reconstruct at least a single genotype for every donor in the selected pool.

A prototypical algorithmic implementation was built using the open source Posgresql database system. The convergence criterion $\epsilon$ was set to $3.8 \times 10^{-4}$% of the final log-likelihood value. Moreover, a threshold, below which a probability was taken to equal zero, was set to $1 \times 10^{-8}$. A naive implementation of the chEM algorithm converged in 22 iterations in approximately 36 hours. While this execution time seems excessive, the common haplotype model does not need to be re-computed in real time. In addition, much faster convergence is expected when performing only model updates.

The final computed model included 11,458 haplotypes with positive probability. The following is a list of the five most common haplotypes and their probabilities identified for the Czech population.

<table>
<thead>
<tr>
<th>A</th>
<th>B</th>
<th>C</th>
<th>DRB1</th>
<th>DQB1</th>
<th>Prob.</th>
</tr>
</thead>
<tbody>
<tr>
<td>01:01</td>
<td>08:01</td>
<td>07:01</td>
<td>03:01</td>
<td>02:01</td>
<td>0.050</td>
</tr>
<tr>
<td>03:01</td>
<td>07:02</td>
<td>07:02</td>
<td>15:01</td>
<td>06:02</td>
<td>0.012</td>
</tr>
<tr>
<td>02:01</td>
<td>13:02</td>
<td>06:02</td>
<td>07:01</td>
<td>02:01</td>
<td>0.010</td>
</tr>
<tr>
<td>23:01</td>
<td>44:03</td>
<td>04:01</td>
<td>07:01</td>
<td>02:01</td>
<td>0.009</td>
</tr>
<tr>
<td>01:01</td>
<td>57:01</td>
<td>06:02</td>
<td>07:01</td>
<td>03:03</td>
<td>0.009</td>
</tr>
</tbody>
</table>

### D. Validation

The method was tested on real patient case studies and the matching results were compared against outputs of existing intermediate level search tools. In searching for matches, a larger donor pool was used. Any donor typed for at least HLA-A/B/DRB1 was considered in the search (this included 90% of the donors contained in the CNRMD database). The ignored database entries include 1,595 donors typed serologically only for HLA-A/B and 3,025 donors typed for HLA-A/B/C, where often the HLA-C typing is erroneous. In practice, the minimal match criterion used is 5 out of 6 matching alleles. Hence, the ignored donors are most likely not suitable candidates. Of the donors that were included in the search, 3,168 were typed at high resolution, 22,848 were typed for HLA-A/B/C/DRB1, and 13,299 were typed for HLA-A/B/DRB1.

The search method currently in use was implemented over 15 years ago to find matching donors using their serological typing data. The method relies on uniform typing data in terms of resolution and HLA loci. Therefore, typing of donors obtained using molecular methods is converted to serological format using standard conversion tables and genotypes are only matched across HLA-A/B/DRB1 loci. A donor with potentially matching alleles at all three loci and both chromosomal copies is called a 6/6 match.

The chEM method presented herein searches for matches across all loci, doesn’t require uniform typing data, and is capable of finding high-resolution matches. For patients with rare genotypes, the method simply ranks the potential donors. For patients with common genotypes, the method ranks the donors and lists their match probabilities. Nonetheless, for comparison purposes, the method was naturally reduced to search for matches across the HLA-A/B/DRB1 loci.
The CNMRD database 6/6 matches identified by the two methods were compared for 100 patients chosen at random. The existing method relies on serological conversion. Hence it is expected that the presented chEM method will shrink the set of potential donors. The table below summarises this comparison study.

<table>
<thead>
<tr>
<th>#</th>
<th>RESULT</th>
</tr>
</thead>
<tbody>
<tr>
<td>41</td>
<td>cases where matches were found using existing method</td>
</tr>
<tr>
<td>29</td>
<td>cases where matches were found using the chEM method</td>
</tr>
<tr>
<td>21</td>
<td>cases where chEM decreased the number of matching donors</td>
</tr>
</tbody>
</table>

As expected, in the majority of cases, the chEM method reduced the set of matching donors. In no cases, the chEM method identified matching donors not identified by the existing method. Out of the cases where the chEM method decreased the number of potential matches, the average reduction was approximately 25.8%.

Below is an example patient and the most likely 6/6 matching donor identified by the chEM method. Despite the donor being typed only at intermediate resolution (the donor has $2 \times 10^{20}$ possible genotypes) for the three loci HLA-A/B/DRB1, the probability of a 6/6 match is equal to 0.97. Using the chEM method, a search for 10/10 matches was also performed. The same donor was also the most likely 10/10 match with match probability equal to 0.11.

<table>
<thead>
<tr>
<th>A</th>
<th>B</th>
<th>C</th>
<th>DRB1</th>
<th>DQB1</th>
</tr>
</thead>
<tbody>
<tr>
<td>02:01</td>
<td>27:05</td>
<td>01:02</td>
<td>01:01</td>
<td>05:01</td>
</tr>
<tr>
<td>02:01</td>
<td>38:01</td>
<td>12:03</td>
<td>13:01</td>
<td>06:03</td>
</tr>
</tbody>
</table>

<table>
<thead>
<tr>
<th>A</th>
<th>B</th>
<th>C</th>
</tr>
</thead>
<tbody>
<tr>
<td>02:</td>
<td>27:</td>
<td>13:</td>
</tr>
<tr>
<td>KDMZ</td>
<td>KCXM</td>
<td>KEGB</td>
</tr>
<tr>
<td>02:</td>
<td>38:</td>
<td>01:</td>
</tr>
<tr>
<td>KDMZ</td>
<td>KCXZ</td>
<td>KEFC</td>
</tr>
</tbody>
</table>

### III. Discussion

One of the fundamental limitations of statistical HLA typing resolution refinement is the computational complexity. Typing ambiguity present in donor pools is almost always prohibitive forcing the use of heuristic reductions. Herein we demonstrate that haplotype models used for refinement need to be only defined over the common haplotypes. As long as the probability of rare haplotypes is sufficiently small, potential donor matches can be identified and sorted purely on the bases of the common haplotype model. A prototypical implementation of the method was proven in terms of robustness and computational speed for the CNMRD database of over 40 thousand donors. The method was further validated in matching donor search case studies, where the results produced using existing donor search tools were significantly refined in the majority of cases.
Denoising of Noisy MRI Brain Image by Using Switching-based Clustering Algorithm

Siti Noraini Sulaiman, Siti Mastura Che Ishak, Iza Sazanita Isa, Norhazimi Hamzah

Abstract— Magnetic Resonance Image is one of technologies used for diagnosing brain cancer. Radiographer use information get from MRI images to diagnose the disease and plan further treatment for the patient. MRI images are always corrupted with noise. Removing noise from images is crucial but it is not an easy task. Filtering algorithm is the most common method used to remove noise. Segmentation technique is normally used to process image in order to detect abnormality happen specifically for brain. But segmentation alone would best to implement when the images is in good condition. In the case where the images are corrupted with noise there are pre processing steps need to be implemented first before we can proceed to the next task. Therefore, in this project we proposed a simpler method that can denoise and at the same time segment the image into several significant region. The proposed method is called switching-based clustering algorithm. The algorithm is implemented on MRI brain images which are corrupted with certain level of salt-and-pepper noise. During the segmentation process, the results showed that the proposed algorithm has the ability to minimize the effect of noise without degrading the original images. The density of noise in the MRI images made varies from 5% to 20%. The results are compared with conventional clustering algorithm. Based on the experimental result obtained, switching-based algorithm provided a better segmentation performance with fewer noise effects than conventional clustering algorithm. Quantitative and qualitative analysis showed positive results for the proposed switching-based clustering algorithm.

Keywords— Magnetic Resonance Image (MRI), Image processing, Image Segmentation, Switching-based Clustering Algorithm, Salt-and-Pepper Noise.

I. INTRODUCTION

Medical imaging is an important instrument in medicine. Magnetic Resonance Image (MRI) is an imaging method commonly used in medical setting to generate high quality images and provide more effective information of the inside of the human body[1, 2]. It presents the clinician with a number of corresponding quick, precise and flexible diagnostic tools. Magnetic Resonance Image is supposed to be very potential for precise measurement of organ anatomy in a simply way. It is very important to obtain correct image in order to make easy the accurate observations for a given application in medical image processing [3]. Besides that, it can also be used to determine the abnormal and normal types of brain.

Unfortunately, medical images are prone to be corrupted by noises. Noise is the unwanted signal present in original signal. It generated in any form or any time. A noisy image can be caused by a faulty memory, a faulty sensor in camera and noise in a channel during transmission. In medical images, salt-and-pepper noise is a typically noise seen on images. It is randomly occurring black or white (or both) pixels over the images. This degradation can be caused by errors in data transmissions and disturbances in the images [4]. Low quality of image is a problem for efficient feature extraction, analysis, and quantitative measurements. Removing noise from medical image is the major task and to denoise the image is not easy. Denoising should be used to improve the image quality for more precise diagnosis [5, 6].

Denoising means minimize the noise in homogenous area without degrading the image details. Denoising is commonly used for post processing method like restoration, segmentation, classification, pattern analysis and others [7]. Conventionally, noise images have to undergo pre-processing step before being subjected further specifies to analysis. Filtering algorithm is the most common method used to remove the noise as pre-processing. In medical imaging, image filtering algorithm technique is used to enhance image quality, increase visibility which helps in the diagnostic process [5]. However those methods are only for filtering for noise removal and applied the further task on the images as image segmentation [8]. Denoising based segmentation method is proposed to be used to remove the noise and at the same time segment the image into several significant regions.

Segmentation algorithm is the method of subdivides an image into different regions [9]. It is widely used in the field of biomedical application [1]. The image segmentation is mainly used to simplify or change an image into something that is more meaningful and easier to analyse without degrading the fine details image s[2]. The radiologists used this method to segment the input medical image meaningful regions. Clustering is the process classification design or pattern of
group which is of similar characteristic for further analysis [10] and widely used for segmentation of MRI brain images. Besides that, clustering algorithm are planned with the objective to create better segmentation [11]. Clustering algorithm figure the typical each cluster in image segmentation. Then, segments the image by categorize each pixel in the closet cluster according to a distance metric. A good results segmentation can be achieved through clustering technique [9]. In MRI brain images, segmentation technique is performed to detect abnormality in brain. It is used to identify three normal tissue such as GM (Gray Matter), WM (White Matter), and CSF (Cerebrospinal Fluid).

In this paper, a new method to minimize noise in the corrupted MRI brain images is proposed i.e. by using a clustering algorithm. The method is known as switching based clustering algorithm [8]. Salt-and-pepper is the fix value impulse noise. It can severely damage data embedded or the information in the original image. For conventional method, when salt-and-pepper noise is occur on the MRI brain images, segmentation can only be apply after complete pre-processing tasks such as algorithm. New clustering algorithm for the segmentation of noisy images is proposed to tackle this problem. The switching-based clustering algorithm allows the clustering algorithm to segment the noise image of low-density salt-and-pepper noise without any prior filtering stage. The difficulty of the process can be reduced by include two independent process (i.e. noise filter and segmentation into one method [8, 9].

This paper organized as follow. In section II explain the switching-based clustering algorithm. Section III describes the proposed methodology. Section IV explains the experiments results and analysis and section V conclusion of paper is discussed.

II. SWITCHING-BASED CLUSTERING ALGORITHM

New method of clustering based segmentation method is proposed for images corrupted with different level salt-and-pepper noise. The method known as switching-based clustering algorithm [8] which is proposed to minimize salt-and-pepper noise sensitivity in segmentation, at the same time increase the potency of the clustering with respect to noise. The algorithm is chosen due to its ability to minimize the noise during the segmentation process without corrupting the fine details of image. The method is divided into two stages. The first stage is the noise detection stage to scan the location of salt-and-pepper noise and record its intensity. The second stage is the clustering process.

In the most clustering algorithm, the algorithms are very sensitive to noise where no spatial information of the image is considered. To minimize the noise sensitivity of the clustering algorithm, the spatial information of the image is commonly considered. This concept is implemented in switching-based AFKM clustering algorithm therefore it is implemented in MRI brain images. The whole implementation of switching AFKM clustering algorithm will be explain in this section as what had stated by [8].

A. Noise Detection

The main purposed is to recognize the intensity of the ‘noise-pixels’ in this stage. Then, place the positions. In the interval [0,255], an 8-bit gray scale digital image is considered with 256 gray levels. Normally, the salt-and-pepper noise has obtained on the high-end and low-end intensities [12-14]. It can be either positive or negative value. The intensity value for the positive impulse is close to 255 (i.e., known as the salt appear white), and the negative impulse with the intensity value of close to 0 (i.e., known as the pepper appears black).

‘Noise pixels’ in the image can be recognizing by these salt-and-pepper noise intensities. According to [13, 14], a binary noise mask N(i, j) will be produced to mark the location of ‘noise-pixels’ by the following:

\[
N(i, j) = \begin{cases} 
0 & \text{if } X(i, j) = L_{\text{salt}} \text{ or } L_{\text{pepper}} \\
1 & \text{otherwise} 
\end{cases}
\]  

(1)

where X(i, j) is the pixel at the location (i, j) with the intensity X, N(i, j) = 1 stand for the ‘noise free pixels’ to be retained in the next clustering stage and N(i, j) = 0 stand for ‘noise pixels’. L_{\text{salt}} and L_{\text{pepper}} are the local maxima representing the highest and lowest intensities.

B. Segmentation via AFKM Clustering

After the noise and noise free pixels is being recognized and identified, the next step i.e. segmentation via clustering is then commence. In this paper, as mention before the AFKM clustering algorithm [15] is chosen to be implemented on the MRI-brain images. The steps for AFKM is explain as follows.

Consider an image with X \times Y pixel as the number columns and Y pixels as the number of rows. Let \( v_t \) be the t-th data and \( c_k \) is the k-th center. The aim of AFKM algorithms is to minimize the objective function given as follows:

\[
J = \sum_{k=1}^{n} \sum_{t=1}^{N} \left( M_{k_t}^{m} \right) \left\| v_t - c_k \right\|^2
\]

(2)

where \left\| \right\| stands for distance measure that is normally considered as the Euclidean norm. All data are assigned to the nearest centre based on the Euclidean distance. Then, m, the fuzziness exponent, is an integer, \( m \geq 1 \). The new position for each centre is calculated using:

\[
c_k = \frac{\sum_{t=1}^{N} \left( M_{k_t}^{m} \right) v_t}{\sum_{t=1}^{N} \left( M_{k_t}^{m} \right)}
\]

(3)
The value of \( M^m_k \) is simplified in the iteration according to [15]:

\[
(M^m_k) = (M^m_k) + \Delta(M^m_k) \tag{4}
\]

As \( \Delta(M^m_k) = \alpha(e_k)(e_k) \) \( \tag{5} \)

where \( \alpha \) is a designed constant with value between 0 and 1 and typically set to 0.1. Then, the value of \( e_k \) is the error Belonging noise. The degree of Belongings, \( e_k \) is the degree of Belongingness, \( B_\alpha \) is calculated by [15]:

\[
e_k = B_k - \hat{B}_k
\]

where \( \hat{B}_k \) is the normalized value for degree of belongingness.

\[
B_k = \left( \frac{c^e}{M^m_k} \right) \tag{6}
\]

The whole implementation of switching-based AFKM can be refer to [8] and [15].

### III. Proposed Methodology

This project is focused on noise removal for medical images (MRI brain) but at the same time segment the images. The implementation of the project is simplified in flow chart below:

#### Figure 1: Methodology

The proposed method to denoise the MRI brain images is by using switching AFKM. In this paper, five of MRI images are chosen from the tested images and to be used to analyse the performance of the proposed algorithm. These images are consists normal and abnormal brain as shown in Figures 2(a) to 2 (e) respectively. The images then are inserted with different level of salt-and-pepper noise (by simulation) using MATLAB 2010 image processing toolbox. The level of noise inserted is between 5% to 20%. After that, the corrupted images is tested with switching-AFKM and compared with the conventional segmentation method i.e the AFKM algorithm. The result are evaluated both qualitative and quantitative analysis.

#### Figure 2: The original image (a) Normal 1, (b) Normal 2, (c) Normal 3, (d) Abnormal 1 and (e) Abnormal 2

### IV. Experimental Results and Analysis

Original MRI brain images were performs on the results as a reference images. The salt-and-pepper noise is added to the images. The corrupted images are 5% to 20% density of salt-and-pepper noise. The noises have to remove to get better results. Denoising is perform by using the proposed algorithm (switching-based clustering algorithm) and compare with conventional clustering algorithm. Comparison being made based on the qualitative and quantitative analysis.

#### A. Qualitative Analysis

Segmented images are evaluated in this analysis. This analysis is depends on the type of purpose on the MRI brain images. The results of qualitative will be different for different application on the images. This is because each method has the good quality and performance of segmentation.

In experiments, the proposed switching-based clustering algorithm on normal and abnormal MRI brain images are corrupted with different levels densities of salt-and-pepper noise. Its results were compared with conventional-based clustering algorithm. The main goal is to determine the performance of the proposed algorithm should minimize the effect of salt-and-pepper noise during segmentation process. All images are segmented into three clusters.

The segmentation results for Normal 1, Normal 2, Normal 3, Abnormal 1 and Abnormal 2 images are shown in Figures 3, 4, 5, 6, and 7 respectively. Image (a) represents the original
image, image (b) is a noisy image, image (c) and (d) a resultant images after applying the proposed switching AFKM and conventional algorithm. The images show different result of performance for different algorithm.

![Segmentation results](image)

Figure 3: Segmentation results of Normal 1 MRI brain image with 5% to 20% density of salt and pepper noise using: (a) original images, (b) noisy images, (c) conventional AFKM and (d) switching AFKM.

Figures 3 (a), 4 (a) and 5 (a) shows the Normal 1, 2 and 3 MRI brain image. Noisy image is show in Figures 3 (b), 4 (b) and 5 (c). The images corrupted with 5%, 10%, 15% and 20% density of salt-and-pepper noise. The images became highly corrupted when noise level is increase. Their results are compared with the proposed switching-AFKM clustering algorithm and conventional-based clustering algorithm. When the image is processed by conventional method in Figures 3 (c), 4(C) and 5 (d), the image is not clear and cannot segment the image in three regions which is (GM) gray matter, (WM) white matter and CSF (Cerebrospinal Fluid). The proposed switching-AFKM clustering algorithm mostly can minimize the effect of salt-and-pepper noise. For the resultant images in Figures 3(d), 4(d) and 5(d), proposed algorithm produce a better segmentation results with salt-and-pepper noise reduced significantly. It provide sharper and clearer images compare with conventional algorithm.

Let consider the segmentation results on the sample with tumor brain images (i.e., Abnormal 1 and Abnormal 2) as shown in Figures 6 and 7. For the images Abnormal 1 and Abnormal 2, it can be seen the proposed clustering algorithm produced similar result. The effect noisy on MRI brain (i.e., salt-and-pepper noise) can be minimizing by using a suitable algorithm. The resultant images showed clearly when using the proposed clustering method. It can detect the location of brain tumor successfully. The proposed algorithm can segment the image perfectly. It will be easy for radiographer to detect the disease like tumor, stroke, Alzheimer and others. The proposed switching AFKM algorithm is successfully produce a good segmentation performance than conventional.

### B. Quantitative Analysis

The quantitative analysis is applied to the MRI brain images. All images are compared the performance between proposed algorithm and conventional clustering algorithm. The evaluation of quantitative is measure between two method algorithms. The Tables 1, 2, and 3 is present the quantitative measurement for the normal images. Tables 3 and 4 are evaluated for the abnormal images. For Table 6, the average of quantitative is measure.

In this paper, function $F(I)$, proposed by Liu and Yang[16], and $F'(I)$ as well as $Q(I)$, proposed by[17] used for evaluation of performances. This function will give result match closely to image decision. Smaller values showed better performance. The three functions are calculated using the equation as follow [8]:

\[
F(I) = \sqrt{R} \sum_{i=1}^{R} \frac{e_i^2}{\sqrt{A_i}}
\]

\[
F(I) = \frac{1}{1000(N \times M)} \left[ \sum_{A=1}^{\text{Max}} \left( \frac{1}{R(A)} \right)^{1+1/A} \sum_{i=1}^{R} e_i^2 \right]
\]

\[
Q(I) = \frac{1}{1000(N \times M)} \sqrt{R} \sum_{i=1}^{R} \frac{e_i^2}{1 + \log A_i + \left( \frac{R(A_i)}{A_i} \right)^2}
\]

Based on three formulas consideration above, I is the resultant image to be calculate, $N \times M$ is the image size, $R$ is the number of region found, $A_i$ is the size of the $i$-th region, and $R(A_i)$ is the number of region having area $A_i$. Furthermore, $e_i$ is defined as the sum of Euclidean distances between the features of pixels of region $I$ and the corresponding region in the segmented image. Based on reference to [18], $R(A)$ is the number of regions having closely area $A$, and Max denotes the area of the largest region in the segmented images.

The proposed AFKM produces the value of $F(I)$, $F'(I)$ and $Q(I)$ are smaller values is better segmentation results. Based on Table 1, the all value of proposed method is smaller than conventional method. This value showed good segmentation results. Referring Tables 2 and 3, the value $F(I)$ and $F'(I)$ on switching AFKM is greater than conventional method when the images is corrupted 20% salt-and-pepper noise and other value is smaller compare both algorithm. However, the qualitative image for proposed method is showed good performance. The qualitative measurement for abnormal image showed in Tables 4 and 5. The parameter proposed by switching AFKM is better than conventional. These support the qualitative results shown in Figure 3 to 7 where the resultant images are showed clearly.
Figure 4: Segmentation results of Normal 2 MRI brain image with 5% to 20% density of salt and pepper noise using: (a) noisy images, (b) conventional AFKM and (c) switching AFKM.

Figure 6: Segmentation results of Abnormal 1 MRI brain image with 5% to 20% density of salt and pepper noise using: (a) noisy images, (b) conventional AFKM and (c) switching AFKM.

Figure 5: Segmentation results of Normal 3 MRI brain image with 5% to 20% density of salt and pepper noise using: (a) noisy images, (b) conventional AFKM and (c) switching AFKM.

Figure 7: Segmentation results of Abnormal 2 MRI brain image with 5% to 20% density of salt and pepper noise using: (a) noisy images, (b) conventional AFKM and (c) switching AFKM.
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Conventional AFKM proposed algorithm. The value of Based on result, the value of Switching Method AFKM showed the smaller when the original is corrupted with 5% and 10% salt-and-pepper noise. Although for $F(I)$ and $F'(I)$ for images corrupted more than 10% is greater, but the qualitative results for images is a better performance. These can be solve by Euclidean distance in term $F(I)$ and $F'(I)$ function. Resulting is smaller if Euclidean distance is smaller. According to Borsotti et al. [17] the $Q(I)$ function are commonly important over the other functions as a guide in regulation segmentation algorithm. The value of $Q(I)$ is lower for almost all images in our proposed algorithm. These can be proving by the qualitative analyses of MRI brain images in section I.

<table>
<thead>
<tr>
<th>Method</th>
<th>Noise</th>
<th>No of Region</th>
<th>$F(I)$</th>
<th>$F'(I)$</th>
<th>$Q(I)$</th>
</tr>
</thead>
<tbody>
<tr>
<td>Conventional AFKM</td>
<td>5%</td>
<td>2651</td>
<td>6292.16</td>
<td>632.563</td>
<td>18726.86</td>
</tr>
<tr>
<td></td>
<td>10%</td>
<td>4674</td>
<td>7885.44</td>
<td>791.536</td>
<td>26827.09</td>
</tr>
<tr>
<td></td>
<td>15%</td>
<td>6302</td>
<td>9198.58</td>
<td>923.052</td>
<td>37401.81</td>
</tr>
<tr>
<td></td>
<td>20%</td>
<td>7251</td>
<td>9235.17</td>
<td>928.817</td>
<td>39538.17</td>
</tr>
</tbody>
</table>

| Switching AFKM | 5% | 211 | 1348.6 | 142.431 | 3099.8 |
|                | 10% | 489 | 3620.31 | 372.831 | 8641.86 |
|                | 15% | 876 | 7386.6 | 752.355 | 18011.76 |
|                | 20% | 1344 | 12344.5 | 1250.64 | 29667.57 |

<table>
<thead>
<tr>
<th>Method</th>
<th>Noise</th>
<th>No of Region</th>
<th>$F(I)$</th>
<th>$F'(I)$</th>
<th>$Q(I)$</th>
</tr>
</thead>
<tbody>
<tr>
<td>Conventional AFKM</td>
<td>5%</td>
<td>4981</td>
<td>14413.33</td>
<td>1446.81</td>
<td>57395.23</td>
</tr>
<tr>
<td></td>
<td>10%</td>
<td>8857</td>
<td>17380.42</td>
<td>1742.62</td>
<td>86695.49</td>
</tr>
<tr>
<td></td>
<td>15%</td>
<td>11754</td>
<td>18481.68</td>
<td>1852.63</td>
<td>114216.3</td>
</tr>
<tr>
<td></td>
<td>20%</td>
<td>13907</td>
<td>18414.13</td>
<td>1845.58</td>
<td>130384.1</td>
</tr>
</tbody>
</table>

| Switching AFKM | 5% | 501 | 3853.86 | 397.588 | 11814.81 |
|                | 10% | 937 | 8428.225 | 859.581 | 26028.9 |
|                | 15% | 1841 | 16037.48 | 1620.83 | 48877.32 |
|                | 20% | 2943 | 26639.19 | 2683.42 | 81046.91 |

<table>
<thead>
<tr>
<th>Method</th>
<th>Noise</th>
<th>No of Region</th>
<th>$F(I)$</th>
<th>$F'(I)$</th>
<th>$Q(I)$</th>
</tr>
</thead>
<tbody>
<tr>
<td>Conventional AFKM</td>
<td>5%</td>
<td>4933</td>
<td>13991.08</td>
<td>1403.78</td>
<td>57344.21</td>
</tr>
<tr>
<td></td>
<td>10%</td>
<td>8732</td>
<td>17133.33</td>
<td>1717.44</td>
<td>87083.1</td>
</tr>
<tr>
<td></td>
<td>15%</td>
<td>11692</td>
<td>18253.96</td>
<td>1829.18</td>
<td>117286.1</td>
</tr>
<tr>
<td></td>
<td>20%</td>
<td>13875</td>
<td>17429.9</td>
<td>1746.8</td>
<td>129781.2</td>
</tr>
</tbody>
</table>

| Switching AFKM | 5% | 352 | 3042.501 | 317.14 | 9590.247 |
|                | 10% | 844 | 7759.901 | 791.561 | 24423.39 |
|                | 15% | 1711 | 15631.24 | 1580.63 | 49534.54 |
|                | 20% | 2860 | 27197.5 | 2739.98 | 86348.29 |

<table>
<thead>
<tr>
<th>Method</th>
<th>Noise</th>
<th>No of Region</th>
<th>$F(I)$</th>
<th>$F'(I)$</th>
<th>$Q(I)$</th>
</tr>
</thead>
<tbody>
<tr>
<td>Conventional AFKM</td>
<td>5%</td>
<td>1258</td>
<td>10362.67</td>
<td>1047.48</td>
<td>25914.64</td>
</tr>
<tr>
<td></td>
<td>10%</td>
<td>2240</td>
<td>14102.84</td>
<td>1423.2</td>
<td>35838.29</td>
</tr>
<tr>
<td></td>
<td>15%</td>
<td>2952</td>
<td>22153.52</td>
<td>2230.11</td>
<td>73907.51</td>
</tr>
<tr>
<td></td>
<td>20%</td>
<td>3590</td>
<td>17883.86</td>
<td>1800.08</td>
<td>46201.2</td>
</tr>
</tbody>
</table>

| Switching AFKM | 5% | 148 | 962.4012 | 102.522 | 2226.688 |
|                | 10% | 381 | 2614.168 | 269.718 | 6149.047 |
|                | 15% | 694 | 5643.765 | 575.518 | 13417.94 |
|                | 20% | 1150 | 9974.313 | 1013.49 | 23986.51 |

Lastly based on the table 6, the average quantitative analysis is calculated. It is to support the qualitative measurement. It tested 28 images which contains normal and abnormal images. Based on result, the value of $Q(I)$ parameters is lower by the proposed algorithm. The value of $F(I)$ and $F'(I)$ parameters showed the smaller when the original is corrupted with 5% and 10% salt-and-pepper noise. Although for $F(I)$ and $F'(I)$ for images corrupted more than 10% is greater, but the qualitative results for images is a better performance. These can be solve by Euclidean distance in term $F(I)$ and $F'(I)$ function. Resulting is smaller if Euclidean distance is smaller. According to Borsotti et al. [17] the $Q(I)$ function are commonly important over the other functions as a guide in regulation segmentation algorithm. The value of $Q(I)$ is lower for almost all images in our proposed algorithm. These can be proving by the qualitative analyses of MRI brain images in section I.

<table>
<thead>
<tr>
<th>Method</th>
<th>Noise</th>
<th>No of Region</th>
<th>$F(I)$</th>
<th>$F'(I)$</th>
<th>$Q(I)$</th>
</tr>
</thead>
<tbody>
<tr>
<td>Conventional AFKM</td>
<td>5%</td>
<td>5330.46</td>
<td>7125.115</td>
<td>716.39</td>
<td>27500.9</td>
</tr>
<tr>
<td></td>
<td>10%</td>
<td>5877.21</td>
<td>10097.43</td>
<td>1014.35</td>
<td>49264</td>
</tr>
<tr>
<td></td>
<td>15%</td>
<td>7825.35</td>
<td>11312.64</td>
<td>1135.75</td>
<td>67740.1</td>
</tr>
<tr>
<td></td>
<td>20%</td>
<td>9223.37</td>
<td>11468.92</td>
<td>1151.33</td>
<td>76860.3</td>
</tr>
</tbody>
</table>

| Switching AFKM | 5% | 483.31 | 2942.296 | 307.77 | 8919.55 |
|                | 10% | 834.71 | 6460.74 | 664.493 | 20501.9 |
|                | 15% | 1412 | 13130.27 | 1979.52 | 3857.83 |
|                | 20% | 2131.43 | 20174.44 | 2045.94 | 66270.4 |

Note: Numbers in bolded font represent the best result.

V. CONCLUSION

This paper presents a proposed algorithm named the switching AFKM for the segmentation process and conventional clustering algorithm. These techniques are applied on normal and abnormal MRI brain image is corrupted with salt-and-pepper noise. The proposed method produced better results in a qualitative analysis. It is able to minimize the salt-and-pepper noise. The image is clear and sharpness than other method. Besides, the good performance of proposed method is proven by quantitative measurement.

REFERENCES


**BIOGRAPHIES**

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Energy Efficiency of Conservative Tillage Systems in the Hilly Areas of Romania

Rusu Teodor

Abstract—Conservative tillage systems, specific to sustainable agriculture, require productivity at least equal to that of conventional technology, optimized energy efficiency and, at the same time, diminished environmental impact. An energy saving way is that of implementing optimal technology specific to each culture and pedoclimatic area. The minimum tillage and no-tillage systems represent alternatives to the conventional system of soil tillage, due to their conservation effects on soil features and to the assured productions, maize: 96-98.1% at minimum tillage and 99.8% at no-tillage, soybean: 102.9-111.9% at minimum tillage and 117.2% at no-tillage, wheat: 93.4-96.8% at minimum tillage and 106.9% at no-tillage, as compared to the conventional system. Correct choice of the right soil tillage system for the crops in rotation help reduce energy consumption, thus maize: 97.3-97.9% at minimum tillage and 91.3% at no-tillage, soybean: 98.6-98.2% at minimum tillage and 92.8% at no-tillage, wheat: 97.4-98% at minimum tillage and 91.6% at no-tillage. Energy efficiency is in relation to reductions in energy savings, but also with efficiency and impact on the tillage system on the cultivated plant. For all crops in rotation, energy efficiency (energy produced from 1 MJ consumed) was the best in no-tillage and 10.44 MJ ha⁻¹ at maize, 6.49 MJ ha⁻¹ at soybean, 5.66 MJ ha⁻¹ at wheat. Energy-efficient agricultural system: the energy consumed-energy produced-energy yield, necessarily have to be supplemented by soil energy efficiency, with the conservative effect of the agricultural system. Only then the agricultural system will be sustainable, durable in agronomic, economic and ecological terms. The implementation of minimum and no-tillage soil systems have increased the organic matter content from 2 to 7.6% and water stable aggregate content from 5.6 to 9.6%, at 0-30 cm depth, as compared to the conventional system. While the soil fertility and the wet aggregate stability have initially been low, the effect of conservation practices on the soil characteristics led to a positive impact on the water permeability in the soil. Availability of soil moisture during the crop growth led to a better plant watering condition. Subsequent release of conserved soil water regulated proper plant water condition and soil structure.

Keywords—No-tillage, minimum tillage, yield, energy efficiency, soil conservation.

I. INTRODUCTION

Sustainable agricultural activity must be organized in a system, scheduled in a sequence and always analysed as part of the relationship: soil-plant-climate area-socio-economic conditions-crop-efficiency [1], [4], [6], [24]. Recommendation of flexible and multifunctional technologies consequently equally aims at reducing the consumption of energy, particularly in the field of aggressive soil tillage, as well as obtaining high yields, soil conservation and environmental protection [2], [8], [11], [13].

The essence of the living system consists in the unique capacity of plants to convert, through photosynthesis, the solar energy, carbon dioxide and water into biochemical alimentary energy. Therefore, a successful measure in agriculture is the quantity of energy gathered under the form of biomass, as a result of efficient human and fossil energy use [3], [5], [9], [10], [12].

The soil tillage has as main purpose a series of immediate effects (with a positive side), resulting from the objectives of the soil tillage themselves: basic tillage, germinal layer preparation, field maintenance. Still, the effects of the soil tillage can often have an immediate negative part or long lasting effects, positive or negative [13], [15], [16], [19], [23], [26], [30].

The influence of soil tillage system on soil properties and energy efficiency is proved by important factors of soil fertility conservation and evaluation of the sustainability of agricultural system [21], [25], [28], [29]. Long-term field experiments provide excellent opportunities to quantify the long-term effects of soil tillage systems on accumulated soil water [18], [20], [22]. The hydrological function of the soil (especially the capacity to retain optimum water quantity, and then gradually make this available for plant consumption) is one of the most important functions determining soil fertility, productivity and soil evolution. Intrinsic soil properties such as organic matter and texture, along with applied tillage practices combine to modify the soil structure, porosity, permeability and water capacity. This, in turn, is a critical factor in the water cycle and affects water accumulation in the soil. The conservation of soil fertility requires a tillage system that optimizes the plant needs in accordance with the soil modifications, that ensures the improvement of soil features and obtaining large and constant crops. Thus, the conservation of soil fertility is tied to maintaining and improving the soil fertility indices and to the productivity of the tillage system.

II. MATERIALS AND METHODS

The experiments have been conducted at the University of Agricultural Sciences and Veterinary Medicine in Cluj Napoca, Romania (46º46'N, 26º36'E), on a moderately fertile Fluvisol [32]. The humus content was 3.01%, pH was 7.2, and soil texture was clay (42% clay in the arable stratum). The
experimental field has an annual temperature of 8.2°C and annual rainfall of 613 mm.

Treatments used in the study were as follows: A. Conventional tillage (CT): V₁-classic plough (20-25 cm) + disc harrow-2x (8 cm) (witness treatment). B. Minimum tillage (MT): V₂-paraplow (18-22 cm) + rotary harrow (8 cm); V₂-chisel plough (18-22 cm) + rotary harrow (8 cm); V₂-rotary harrow (10-12 cm). C. No-tillage (NT): V₃-direct drill with Accord Optima HD for hoeing and SUP adapted for wheat.

All soil tillage was accomplished during the autumn period for wheat; for corn and soybeans we used the plough, paraplow, chisel plough in the autumn and finally, for the germinal layer preparation, we used the disc harrow and rotary harrow in the spring. Crop rotation was: maize-Zea mays L., soy-bean-Glycine hispida L. Merr. and wheat-Triticum aestivum L.

The experimental design was a randomized complete block design with three replications. The area of a parcel was 300 m². Except for the soil tillage system, all other variables were held constant, including the herbicide used: wheat-post emergent dicamba 120 g/l + 2.4D 300 g/l, 0.91 ha⁻¹; corn-pre emergent acetochlor 820-860 g/l + antidote, 2.51 ha⁻¹ and post emergent dicamba 120 g/l + 2.4D 300 g/l, 0.91 ha⁻¹; soybeans-pre emergent acetochlor 820-860 g/l + antidote, 2.51 ha⁻¹ and post emergent bentazon 480 g/l + Wettol 150 g/l, 2.51 ha⁻¹.

To quantify the change in soil properties under different tillage practices, determinations were made for each culture in four vegetative stages (spring, 5-6 leaves, bean forming and harvest). Soil parameters monitored included soil water content (gravimetric method, Aquatek probe-Frequency domain reflectometry), soil bulk density (determined by volumetric ring method using the volume of a ring 100 cm³), water stable aggregates (Czeratzki method), soil permeability (using the Infiltrometer method) and organic matter content (Walkley-Black method). The average values obtained during the vegetal phases were statistically analysed. The results were analysed using ANOVA and Duncan's test [31]. A significance level of P ≤ 0.05 was established a priori.

Regarding energetic assessment, the most realistic means of comparison of various agricultural technologies remains energy efficiency, using the following indicators: Energy Efficiency Factor: e = (Er-Ec)/Er [MJ]; Energy Yield: γ = Er/Ec [MJ]; Energy Report r = Ec/Er [MJ]. Where: Er-energy as gathered biomass [MJ]; Ec-technologically consumed energy to produce this biomass [MJ].

Consumed and produced energy represent in fact a sum of inputs and outputs in the technological process. Consequently: Er = Erp+Ers [MJ]. Where: Erp-energy corresponding to primary harvest; Ers-energy corresponding to secondary harvest.

Technologically consumed energy has several components: Ect=Ect+Ec+Ec+Ec+Ec+Eo [MJ]. Where: Ect-energy consumption related with the tractor [MJ]; Ecm-energy consumption related with agricultural machinery [MJ]; Ec-energy consumption related with seeds [MJ]; Ecf-energy consumption related with fertilization [MJ]; Ecg-energy consumption related with pesticides [MJ]; Ecq-energy consumption related to human work resources [MJ]; Eo-energy consumed in other ways [MJ]. Each component is the sum of elementary energies specific to each technological operation. Quantification of consumed energy and of the produced energy has been achieved on the basis of equivalents mentioned in specialty literature [7], [17], [27].

The equivalence indicators are:


Energy produced-1 kg corn: 16.41 MJ; 1 kg corn cob: 15.29 MJ; 1 kg soy bean: 20.79 MJ; 1 kg soy stems: 15.42 MJ; 1 kg wheat: 16.06 MJ; 1 kg wheat straws: 15.26 MJ.

III. RESULTS AND DISCUSSION

The soil tillage system influences the yields obtained in a differentiated way, depending on the culture type (table 1). Corn crop assures the highest yield with plough and no-tillage systems. Paraplow and chisel give smaller yields (6,710-6,730 kg ha⁻¹), with statistically ensured differences (significantly negative) and confirmed by the test of multiple comparisons, Duncan's test (ab). The smallest corn productions were obtained with rotary harrow, the differences being distinctly negative, statistically ensured (b). Soybean culture had the best reaction within the rotation, both with the no-tillage (very significant positive differences as compared to the plough), as well as with minimum soil tillage system, with paraplow and rotary harrows (ab). For wheat culture no-tillage ensure highest yield, 3,986 kg ha⁻¹, and the lowest production has been achieved with chisel (93.4%).

The quantity of energy produced depending on soil tillage system, is related to main and secondary yield, being higher in the plough variant. Energy efficiency is influenced by the soil tillage system, being higher in no-tillage (e=0.9042, 101%), followed by the variants with chisel and paraplow (100.1%). Energetic efficiency is influenced by the energy consumed within every technologic system, the smaller the consumed energy within the system, the higher the efficiency. The high power efficiency in no-tillage (γ=10.44 MJ ha⁻¹), chisel (γ =9.66 MJ ha⁻¹) and paraplow (γ =9.65 MJ ha⁻¹), as compared to the plough system (γ =9.54 MJ ha⁻¹), shows that the energy invested in these variants has had a higher efficiency. The proportional expression between the produced and consumed energy, through energetic report, as certain the lower value of this indicator of 0.096 in no-tillage and the highest value, of 0.108 for the rotary harrow variant.

Considering the amount of produced energy, in maize culture, we can emphasize the advantages of the plough variant. The intense soil mobilization, in conjunction with the effects produced in the soil linked to the release of adequate
nutrients and providing necessary conditions for maize development ensures the highest productions. Intense impact on soil does not, however, always have positive effects. Eventually, the energy efficiency demonstrates the superiority of the no-tillage and minimum tillage systems, in terms of energy consumption reductions and optimization of agricultural technologic system.

### Table 1. The influence of different soil tillage systems upon the plants yield in the case of maize, soybean and wheat crops.

<table>
<thead>
<tr>
<th>Soil tillage systems</th>
<th>Maize, kg ha$^{-1}$</th>
<th>Significance (%)</th>
<th>Soybean, kg ha$^{-1}$</th>
<th>Significance (%)</th>
<th>Wheat, kg ha$^{-1}$</th>
<th>Significance (%)</th>
</tr>
</thead>
<tbody>
<tr>
<td>Classic plough + disc-2x</td>
<td>6730 ab</td>
<td>98.1</td>
<td>3385 ab</td>
<td>111.9</td>
<td>3615 ab</td>
<td>96.9</td>
</tr>
<tr>
<td>Paraplow + rotary harrow</td>
<td>6730 ab</td>
<td>98.1</td>
<td>3385 ab</td>
<td>111.9</td>
<td>3615 ab</td>
<td>96.9</td>
</tr>
<tr>
<td>Chisel plow + rotary harrow</td>
<td>6730 ab</td>
<td>98.1</td>
<td>3385 ab</td>
<td>111.9</td>
<td>3615 ab</td>
<td>96.9</td>
</tr>
<tr>
<td>Rotary harrow</td>
<td>6730 ab</td>
<td>98.1</td>
<td>3385 ab</td>
<td>111.9</td>
<td>3615 ab</td>
<td>96.9</td>
</tr>
<tr>
<td>No Tillage</td>
<td>6730 ab</td>
<td>98.1</td>
<td>3385 ab</td>
<td>111.9</td>
<td>3615 ab</td>
<td>96.9</td>
</tr>
</tbody>
</table>

Note: wt—witness, ns—not significant, *—positive significance, —negative significance, a, ab, b, c—Duncan’s classification (the same letter within a row indicates that the means are not significantly different)

Maize: DL5% = 100.01 kg ha$^{-1}$, DL1% = 151.45 kg ha$^{-1}$, DL0.1% = 243.30 kg ha$^{-1}$

Soybean: DL5% = 190.75 kg ha$^{-1}$, DL1% = 271.16 kg ha$^{-1}$, DL0.1% = 392.62 kg ha$^{-1}$

Wheat: DL5% = 241.21 kg ha$^{-1}$, DL1% = 338.57 kg ha$^{-1}$, DL0.1% = 477.99 kg ha$^{-1}$

### Table 2. The influence of the soil tillage system on energy efficiency in maize culture.

<table>
<thead>
<tr>
<th>Variant</th>
<th>Energy Efficiency (γ)</th>
<th>Energy Efficiency (%)</th>
<th>Energy Yield (γ)</th>
<th>Energy report (r)</th>
</tr>
</thead>
<tbody>
<tr>
<td>Classic plough + disc-2x</td>
<td>0.8952</td>
<td>100</td>
<td>9.54</td>
<td>0.104</td>
</tr>
<tr>
<td>Paraplow + rotary harrow</td>
<td>0.8963</td>
<td>100</td>
<td>9.65</td>
<td>0.103</td>
</tr>
<tr>
<td>Chisel plow + rotary harrow</td>
<td>0.8965</td>
<td>100</td>
<td>9.66</td>
<td>0.103</td>
</tr>
<tr>
<td>Rotary harrow</td>
<td>0.8915</td>
<td>99.6</td>
<td>9.22</td>
<td>0.108</td>
</tr>
<tr>
<td>No-Tillage</td>
<td>0.9042</td>
<td>101.0</td>
<td>10.44</td>
<td>0.096</td>
</tr>
</tbody>
</table>

The energy required for setting up and maintaining the soybean culture after conventional system represents 25,364.09 MJ/ha and goes down to 97.6-98.2% at NT and at 92.8% at NT. Energy efficiency is superior in all variants as compared to the witness, soy reacting very well with MT and NT systems. Energy yield confirms this positive reaction, the results being 6.49 MJ ha$^{-1}$ at NT and 5.51-5.97 MJ ha$^{-1}$ at MT, for 1 MJ ha$^{-1}$ consumed.

### Table 3. Influence of soil tillage system on energy efficiency in soybean culture.

<table>
<thead>
<tr>
<th>Variant</th>
<th>Energy Efficiency (γ)</th>
<th>Energy Efficiency (%)</th>
<th>Energy Yield (γ)</th>
<th>Energy report (r)</th>
</tr>
</thead>
<tbody>
<tr>
<td>Classic plough + disc-2x</td>
<td>0.8091</td>
<td>100</td>
<td>5.23</td>
<td>0.191</td>
</tr>
<tr>
<td>Paraplow + rotary harrow</td>
<td>0.8325</td>
<td>102.9</td>
<td>5.97</td>
<td>0.167</td>
</tr>
<tr>
<td>Chisel plow + rotary harrow</td>
<td>0.8184</td>
<td>101.1</td>
<td>5.51</td>
<td>0.182</td>
</tr>
<tr>
<td>Rotary harrow</td>
<td>0.8298</td>
<td>102.5</td>
<td>5.88</td>
<td>0.170</td>
</tr>
<tr>
<td>No-Tillage</td>
<td>0.8458</td>
<td>104.5</td>
<td>6.49</td>
<td>0.154</td>
</tr>
</tbody>
</table>

In the case of autumn wheat culture, technology is energetically equivalent to 23,272.38 MJ ha$^{-1}$ through the CT system (table 4). Application of MT reduces energy consumption to 97.4-98%, and NT to 91.6%, compared with the plough system. The influence of the soil tillage system on the amount of gathered energy reflects on the energy efficiency factor, where, in comparison with the witness, a higher efficiency at NT has been calculated (101%). Energy efficiency has been reduced in the other variants, but it does not fall below 99%. Energy yield shows that in 1 MJ ha$^{-1}$ consumed a larger amount of energy is obtained with no-tillage (γ=5.66 MJ ha$^{-1}$), and the lowest yield was recorded with the chisel plough variant, 5.32 MJ ha$^{-1}$. The energy report has the best value in no-tillage (0.177), followed by the plough variant (0.179). Statistical analysis of the results demonstrated that the differences in accumulated soil water depended on the variants of soil tillage (table 5). Soil texture and structure have a strong effect on the available water capacity. The results clearly demonstrate that MT and NT systems promote increased humus content (2-7.6%) and increased water constant aggregate content (5.6-9.6%) at 0-30 cm depth as compared to conventional tillage. Multiple analysis of soil classification and tillage system on the hydric stability of soil structure and water supply accumulated in soil have shown that all variants with minimum tillage are superior (b, c), having a positive influence on soil structure stability. The increase in organic matter content is due to the vegetal remnants at the soil surface (NT) or partially incorporated (MT) and adequate biological activity in this system. In the case of humus content and also in the
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hydro stability structure, the statistical interpretation of the data shows an increasing positive significance of the MT and NT systems application. The soil fertility and wet aggregate stability were initially low, the effect being the conservation of the soil features and also their reconstruction, with a positive influence on the permeability of the soil for water. More aggregated soils permit more water to reach the root zone. This does not only increase productivity, but it also reduces runoff, and thus the erodibility potential.

The bulk density values at 0-30 cm increased by 0.01-0.03% under minimum and no-tillage systems. This raise was not significant in any of the experimental variants. Multiple comparing and classification of experimental variants align all values on the same level of significance (a). On mollic Fluvisol, soils with good permeability, high fertility, and low susceptibility to compaction, accumulated water supply was higher (representing 12.4-15%) for all minimum and no-tillage soil systems.

The minimum tillage and no-tillage systems represent alternatives to the conventional system of soil tillage, due to their conservation effects on soil features and to the assured productions, maize: 96.98-1.1% at MT and 99.8% at NT, soybean: 102.9-111.9% at MT and 117.2% at NT, wheat: 93.4-96.8% at MT and 106.9% at NT, as compared to the conventional system.

Correct choice of the right soil tillage system for the crops in rotation help reduce energy consumption, thus maize: 97.3-97.9% at MT and 91.3% at NT, soybean: 98.6-98.2% at MT and 92.8% at NT, wheat: 97.4-98% at MT and 91.6% at NT. Energy efficiency is in relation to reductions in energy savings, but also with efficiency and impact on the tillage system on the cultivated plant, maize: 99.6-100.1% at MT and 101% at NT, soybean: 101.1-102.9% at MT and 104.5% at NT, wheat: 99-99.9% at MT and 100.4% at NT. For all crops in rotation, energy efficiency (energy produced from 1 MJ consumed) was the best in no-tillage and 10.44 MJ ha\(^{-1}\) at maize, 6.49 MJ ha\(^{-1}\) at soybean, 5.66 MJ ha\(^{-1}\) at wheat.

Energy-efficient agricultural system: the energy consumed-energy produced-energy yield, necessarily have to be supplemented by soil energy efficiency, with the conservative effect of the agricultural system. Only then the agricultural system will be sustainable, durable in agronomic, economic and ecological terms.

IV. CONCLUSIONS

The minimum tillage and no-tillage systems represent alternatives to the conventional system of soil tillage, due to their conservation effects on soil features and to the assured productions, maize: 96.98-1.1% at MT and 99.8% at NT, soybean: 102.9-111.9% at MT and 117.2% at NT, wheat: 93.4-96.8% at MT and 106.9% at NT, as compared to the conventional system.

Correct choice of the right soil tillage system for the crops in rotation help reduce energy consumption, thus maize: 97.3-97.9% at MT and 91.3% at NT, soybean: 98.6-98.2% at MT and 92.8% at NT, wheat: 97.4-98% at MT and 91.6% at NT. Energy efficiency is in relation to reductions in energy savings, but also with efficiency and impact on the tillage system on the cultivated plant, maize: 99.6-100.1% at MT and 101% at NT, soybean: 101.1-102.9% at MT and 104.5% at NT, wheat: 99-99.9% at MT and 100.4% at NT. For all crops in rotation, energy efficiency (energy produced from 1 MJ consumed) was the best in no-tillage and 10.44 MJ ha\(^{-1}\) at maize, 6.49 MJ ha\(^{-1}\) at soybean, 5.66 MJ ha\(^{-1}\) at wheat.

Energy-efficient agricultural system: the energy consumed-energy produced-energy yield, necessarily have to be supplemented by soil energy efficiency, with the conservative effect of the agricultural system. Only then the agricultural system will be sustainable, durable in agronomic, economic and ecological terms.

This study demonstrated that increased organic matter content in soil, aggregation, and permeability are all promoted by minimum and no-tillage systems. The implementation of such practices ensures a greater water supply. The practice of reduced tillage is ideal for enhancing soil fertility, water accumulation capacity, and reducing erosion. The advantages of minimum and no-tillage soil systems for Romanian pedoclimatic conditions can be used to improve methods in low producing soils with reduced structural stability on sloped fields, as well as measures of water and soil conservation on the whole ecosystem.

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A suggested method for assessing cliff instability susceptibility at a given scale (CISA)

G.F. Andriani, V. Pellegrini

Abstract—This paper illustrates a new multidirectional method for assessing cliff instability susceptibility at a given scale (CISA, Cliff Instability Susceptibility Assessment) through a case study along the Murgia coastline North of Bari (Apulia, SE Italy). The stretch of coastline considered in this study shows an indented rocky coast with cliffs up to 12 m high, numerous small to medium-sized caps and inlets, and a well protected tourist port. The coastal outcrops are made up of Mesozoic carbonate rocks which are thickly-bedded (0.2-1.0 m) and, moderately to highly fractured and karstified. At places, clear signs of coastal erosion are evident; they consist mainly of rock falls caused by differential erosion of rock strata of varying resistance to weathering and sea wave action on the cliff face.

The CISA method is based on classifying coastal sectors using above all morphological criteria and characterising them estimating and combining 28 incidence parameters according to an heuristic approach. These parameters were divided in four categories: geomechanical (12), morphological (6), meteo-marine (8) and anthropogenic (2). For each parameter 5 classes of rating were proposed; the cliff classification, in terms of cliff instability susceptibility, was obtain from the total rating which represents the summation of the single rating of the individual parameter.

Keywords—Cliff, Carbonate, Instability, Method.

I. INTRODUCTION

In many sites of Apulia (SE Italy), the coastline is gradually receding inland as a result of natural and human processes.

In particular, morphologic features of the rocky coast, mostly produced by Quaternary tectonics, are strongly conditioned by complex mechanisms involving sea waves action against the cliffs, carbonation, weathering and urbanization pressure [1]. The dominant and more visible retreat process of the cliffs consists of slope mass movements of different types and sizes, which include rockfalls, topples and slides controlling by discontinuity pattern and density, and mechanical properties of the carbonate outcrops. Actually, many are the basic factors which play an important role in the assessment of the retreat mechanisms of the Apulian rocky cliffs and the mapping of the coastal stretches most susceptible to erosion is not a simple matter. Different methods are adopted for determining potentially instable areas or landslide hazard assessments; these techniques can be divided into three groups: expert evaluation, statistical methods, and mechanical approach [2]-[3]-[4]-[5] and references therein.

II. SETTING

The study area is located on the Adriatic side of the Murge plateau, an emerged part of the Apulian foreland, about 20 km NW of Bari, in the territory of the Municipality of Giovinazzo (Fig. 1). The Murge plateau is characterized by a 3 km thick
Cretaceous shallow-water succession of limestones and dolostones forming a S-SW dipping monocline slightly deformed by folds and sub-vertical normal and transtensional faults [6]. This succession consists of lagoonal and peritidal carbonates mostly deposited in low-energy inner-platform environments [7]. Only few stratigraphic intervals contain more open- and deeper-marine lithofacies in which abundant and diversified associations of rudists and benthic foraminifera can be found [8]. The limestones and dolostones are bedded, jointed and subjected to karst processes, and represent a peculiar hydrogeological domain. In fact, the hydraulic base level of groundwater circulation corresponds to sea level. Several coastal springs drain groundwater along preferential pathflows where rock-mass permeability is greater [1]. The prevailing morphologic characteristic of the Apulian coastal area is the presence of a series of marine terraces linked by small scarps subparallel to the coastline. These are carved by short erosive incisions (locally named “lame” and “gravine”) in simple catchments and watersheds that are difficult to recognize, as often occurs in karst areas [9]-[10]. The marine terraces develop on wide plains from about 150 m a.s.l. to the present sea level, maintaining a gentle slope to the NE, and give the southeastern side of the Murge a typical terraced profile.

The coastal stretch of Giovinazzo is characterized by the outcropping of the lower part of the Calcare di Bari Fm. (Callovian pp.-early Turonian). This part consists of an about 25-m-thick interval characterized by rudist-dominated medium- to coarse-grained deposits. Rudists form dm- to m-thick sheetlike tabular bodies with whole shells in growth position or more commonly randomly oriented [7]-[8]. The rudist beds gradually pass upward to mud-dominated fine-grained biopeloidal mudstones/wackestones and to laminated fenestral peloidal bindstone showing microbial laminations (Fig. 2).

Morphological features along the Giovinazzo shoreline are typical of an indented rocky coastline where caps and inlets follow each other and steep cliffs, here from about 1.0 m (microcliffs) to 10 m in height, end in subhorizontal surfaces at the top and bottom. Before the tourist port, further to NW, the cliff is protected by a subhorizontal or gently sloping inward wave-cut platform; the last extends from the base of the cliff for 5 m to 10 m and is submerged or emerged with rising and falling of the sea level. The most part of the coastal stretch is protected by a retaining wall made of carbonate stone masonry or unreinforced concrete (Fig. 3).

The wall is about 12 m at its highest part where there is the south side of the historical center of the city which lies on a promontory overlooking the sea. The north side of the historical center overlooks the tourist port which is bordered seaward by two jetties. Large rock or concrete boulders chaotically arranged were used in the past for the protection of the cliff and the outer walls of the old city (Fig. 4).

To the SE, a bay characterized the whole sector that at
places presents a small beach mainly formed by calcareous
and concrete pebbles coming from coastal erosion processes
and improvised attempts of beach nourishment (Fig. 5).

III. METHOD

The CISA method (Cliff Instability Susceptibility Assessment) is a multidirectional method for assessing cliff
instability at a given scale. The first step of the method
consists of the subdivision of the coastal stretch in coastal
sectors based above all on morphological criteria. Cliff height,
coastal landforms (caps, inlet etc.), wave and wind exposure,
coastal defence structures, natural and artificial in types
(wave-cut platform, beach, rock blocks, jetties etc.),
took into account to subdivide in 5 sectors the coastal stretch
under consideration. “Regional Technical Map” (CTR) of
Apulia at scale 1:5000 (www.sit.puglia.it), georeferenced to
WGS84/UTM zone 33N system, was used as base map to
generate the Coastal Stability Map of the study area (Fig. 6).

In the second phase, a qualitative assessment of the cliff
instability integrating traditional geomechanical surveys was
completed by geotechnical laboratory tests, deterministic
analysis (the estimate of the critical height for vertical slope
using the lower bound theorem of limit analysis), multifactor
spatial GIS analysis using physical geographically-
based measures with the purpose of assigning the right weight
to the parameters considered in this study. 28 incidence
parameters were considered and regrouped in four categories:
geomechanical (12), morphological (6), meteo-marine (8) and
anthropogenic (2). For each parameter 5 classes of rating were
proposed; the cliff classification, in terms of instability
susceptibility, was obtain from the total rating which
represents the summation of the single rating of the individual
parameter (Table I). The stability classes with respect to
coastal erosion are reported in Table II.

Discontinuities in rock masses were described according to
the ISRM standards [11]. Favorability/unfavorability of
discontinuities related to cliff stability was evaluated on the
basis of the ratio between discontinuity orientation and
persistance, and potential failure mechanism. With respect to
parallelism between joints and slope face strikes, the presence
of tension cracks at the cliff-top edge was considered the most
hazardous condition. The RQD estimation was carried out by
the volumetric joint count (Jv) and block sizes [12].
Notches and karst features were account as prominent
natural coastal hazards.

Following the standard test procedure outlined in ISRM [13]-[14],
dry unit weight (γd), porosity (n) and uniaxial compressive strength in the dry state (σc) were
determined on 10 cylindrical specimens (100 mm in diameter)
prepared from little rock blocks fallen from the cliff face and
collected along the coastline. As regards the specific gravity
(G), reference was made to a value of 2.70 on the basis of the

Table I – Rating of 28 parameters of the CISA method.

<table>
<thead>
<tr>
<th>Parameter</th>
<th>Very bad</th>
<th>Bad</th>
<th>Normal</th>
<th>Good</th>
<th>Very good</th>
</tr>
</thead>
<tbody>
<tr>
<td>Geomechanical</td>
<td></td>
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<td></td>
<td></td>
</tr>
<tr>
<td>Number of joint</td>
<td></td>
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<td></td>
<td></td>
</tr>
<tr>
<td>Joint spacing</td>
<td></td>
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<td></td>
<td></td>
<td></td>
</tr>
<tr>
<td>Joint surface</td>
<td></td>
<td></td>
<td></td>
<td></td>
<td></td>
</tr>
<tr>
<td>Joint persistence</td>
<td></td>
<td></td>
<td></td>
<td></td>
<td></td>
</tr>
<tr>
<td>Joint roughness</td>
<td></td>
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<tr>
<td>Joint water-filled</td>
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<tr>
<td>Joint weathered</td>
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</tr>
<tr>
<td>Joint oriented to the cliff</td>
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</tr>
<tr>
<td>Joint oriented to the slope face</td>
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<td></td>
<td></td>
<td></td>
</tr>
<tr>
<td>Joint oriented to the sea</td>
<td></td>
<td></td>
<td></td>
<td></td>
<td></td>
</tr>
<tr>
<td>Joint oriented to the wind</td>
<td></td>
<td></td>
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<td></td>
</tr>
</tbody>
</table>

Table II – Stability classes as per CISA values.

<table>
<thead>
<tr>
<th>CISA Value</th>
<th>Stability</th>
</tr>
</thead>
<tbody>
<tr>
<td>&lt; 30</td>
<td>Totally unstable</td>
</tr>
<tr>
<td>30-60</td>
<td>Unstable</td>
</tr>
<tr>
<td>60-90</td>
<td>Stable</td>
</tr>
<tr>
<td>90-120</td>
<td>Very stable</td>
</tr>
<tr>
<td>120-140</td>
<td>Completely stable</td>
</tr>
</tbody>
</table>

Discontinuities in rock masses were described according to
the ISRM standards [11]. Favorability/unfavorability of
discontinuities related to cliff stability was evaluated on the
basis of the ratio between discontinuity orientation and
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of tension cracks at the cliff-top edge was considered the most
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Notches and karst features were account as prominent natural
coastal hazards. Following the standard test procedure
outlined in ISRM [13]-[14], dry unit weight (γd), porosity (n)
and uniaxial compressive strength in the dry state (σc) were
determined on 10 cylindrical specimens (100 mm in diameter)
prepared from little rock blocks fallen from the cliff face and
collected along the coastline. As regards the specific gravity
(G), reference was made to a value of 2.70 on the basis of the
chemical composition of the rocks. These last, in fact, are composed mostly of carbonate with an negligible insoluble residue [15]. Dry unit weight is in the range 19.0-23.2 kN/m³, porosity between 12.2-18.1% and uniaxial compressive strength between 18.3-112.5 MPa. Rock mass hardness was measured in situ with the Schmidt hammer (type L). Higher porosity and, lower strength and density values are due to high weathering rates of a coarse-grained bed at the base of the cliff cropping out in the first coastal sector. It should be noted that in this paper the term weathering includes carbonation. For each coastal sector, the estimate of the critical height of the cliff was carried out with the lower bound theorem of limit analysis adopting the Mohr-Coulomb failure criterion. The shear strength of the rock masses was obtained with the envelope derived by [16]. A precautionary approach was adopted in the analysis so that for each lithotechnical unit estimated on the cliff the lowest value of strength was utilised. Two-dimensional cliff stability analysis was performed using the weighted mean for the geotechnical parameters of the different lithotechnical units defined on the cliff face; therefore in the vertical cliff model, shear strength and unit weight were considered as uniformly constant. A correction factor equal to RQD was applied at the weighted mean of the unit weight determined in laboratory on cylindrical samples for assessing the unit weight of the rock mass. For the submerged portion of the cliff face the buoyant unit weight was taken into account. The stability index (Is) was then calculated as the ratio between cliff height and its critical height. The stability index was first proposed by [17], but the method adopted for calculating critical height of cliff face and range of values given for defining stability classes are different from those proposed in this study.

![Fig. 7 – Differential erosion of rock strata of varying resistance to weathering on the cliff face.](image)

Strong side winds and fetch (effective length) were measured for each coastal sector in correspondence of caps, inlets and the tourist port of Giovinazzo. The effective fetch was defined along the NW, N and NE winds using the recommended procedure of the Shore Protection Manual [18] for a mid-latitude semi-enclosed basin such as the Adriatic Sea. The bottom depth and slope were calculated from the bathymetric data determined in GIS environment (Italian Nautical Charts). The maximum of the values of the offshore spectral height, Hs (m), and the offshore time peak, Tp (s), were obtained from the data collected at the Monopoli buoy (Lat 40°58’30”N; Long 17°22’36” E; World Geodetic System 84) of the National Wave Measuring Network (RON) for the period July 1989-April 2008. Finally, the offshore wavelength Lo (m) was obtained from the term Linear (or Airy) Wave Theory [19] along the NW wind (315-0°N), the N wind (315 – 45° N) and the NE wind (0 – 90°N). For the CISA method, the breaking wave depth and the breaking wave height were determined with the Goda's nomographs [20], while the impact wave height was calculated with the empirical relationship developed by [21]; the type of breaking wave was obtained by the Okazaki & Sunamura’s laboratory study [22].

IV. RESULTS AND DISCUSSION

First of all, long-term morphodynamic evolution of the coastal stretch is influenced by human activities because the study area is characterized by a heavily urbanized coastline with seafront buildings and streets. Furthermore, the most part of the coastal stretch studied presents shore protection structures such as jetties and seawalls with natural stone facing (limestones) or unreinforced concrete. Therefore, clear signs of coastal recession are evident and include the breaking down and removal of material along the coastline by the movement of sea-water (sectors IV and V) and rock falls caused by differential erosion of rock strata of varying resistance to weathering and sea wave action on the cliff face (sector I). Data from field observations indicate that the coastal recession appears not uniform with time and mainly governed by cliff collapses. Notches and their increase in size cause rising shear stresses that induce the cliff to fail, but the presence of a highly weatherable laminated bed cropping out along the cliff face is considered a hazardous condition (Fig. 7).

The basic factors controlling the sea cliff recession are the assailing forces of wave and the resisting force of the cliff-forming rock masses. The wave action consists not only of hydraulic actions (compression, tension, cavitation and wear) but also of abrasive action due to wave-moving pebbles and boulders and wedge action due to the air compressed in fissures by waves [23]. Rock mass strength is controlled by discontinuities and mechanical properties of intact rock pieces. Reduction in rock mass strength is due to weathering and fatigue caused by cyclic loading of waves at the cliff base. In particular, in the presence of an emerged wave-cut platform or rock boulders at the cliff base, weathering is the first responsible for the cliff collapses. The weathering processes include carbonation, salt weathering, water layer weathering (associated with the wetting and drying process) and biological weathering, especially by boring organisms [24]-[25]. In this case, the role of meteo-marine parameters appears to be secondary in the cliff erosion and collapses (Fig. 8).

The assailing forces of wave depend on the wave energy, in turn depending on wind strength and duration, water depth and density, and fetch. The intensity of erosive forces controlling failure mechanisms is determined by the wave type immediately in front of the cliffs and this is determined by the
relationship among offshore wave characteristics (wave height, wave angle and wave period), tidal condition and nearshore submerged morphology.

As a result of the analysis of the data obtained from the applications of the CISA method, it was found that the coastal stretch of Giovinazzo is unstable, although at places retreat processes are opposed by engineering protection structures (jetties, retaining walls etc.). Furthermore, the urbanization of the coastal area has de facto prevented the beach accretion at the cliff toe as a natural barrier against wave action and shoreline erosion.

Table III – Partial rating and CISA rating obtained for the study area.

<table>
<thead>
<tr>
<th>Parameters</th>
<th>Coastal sectors</th>
<th>I</th>
<th>II</th>
<th>III</th>
<th>IV</th>
<th>V</th>
</tr>
</thead>
<tbody>
<tr>
<td>Over-arched</td>
<td>2</td>
<td>3</td>
<td>2</td>
<td>3</td>
<td>2</td>
<td></td>
</tr>
<tr>
<td>Spacing</td>
<td>2</td>
<td>2</td>
<td>2</td>
<td>3</td>
<td>1</td>
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</tr>
<tr>
<td>Severe</td>
<td>2</td>
<td>3</td>
<td>3</td>
<td>3</td>
<td>3</td>
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</tr>
<tr>
<td>Roughtess</td>
<td>3</td>
<td>3</td>
<td>3</td>
<td>3</td>
<td>3</td>
<td></td>
</tr>
<tr>
<td>Wave angle</td>
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<td>4</td>
<td>3</td>
<td>3</td>
<td>3</td>
<td></td>
</tr>
<tr>
<td>Geology</td>
<td>1</td>
<td>3</td>
<td>1</td>
<td>3</td>
<td>2</td>
<td></td>
</tr>
<tr>
<td>Precipitation between joints and slope face cracks</td>
<td>1</td>
<td>3</td>
<td>3</td>
<td>3</td>
<td>2</td>
<td></td>
</tr>
<tr>
<td>Joint weathering with respect to visible local dislocation</td>
<td>1</td>
<td>2</td>
<td>2</td>
<td>2</td>
<td>1</td>
<td></td>
</tr>
</tbody>
</table>

**Table III – Partial rating and CISA rating obtained for the study area.**

**Fig. 8 – Cliff erosion and collapses due to weathering processes on the cliff face.**

In the study area, the influence of the tidal condition was considered significant for the water layer weathering only and this because we are dealing with a typical situation in microtidal environment. Second, applying the CISA method, the weight of the meteo-marine parameters in establishing potential instability appears to be rather uniform in all the sectors with respect to that of the other groups of parameters. In situ observations spanning from 1992 to present allow to affirm that distance of the breaker zone, wave impact height, wave breaker type and exposure to storm wave fronts, expressed in terms of the angle between the coastline and prevailing storm wave fronts, seem to have a different but fundamental incidence on defining the nearshore wave energy for the studied coastal sectors. At the same time, the role of scattering processes induced by nearshore morphology is of great importance in coastal retreat mechanisms, so visual evidence of the wave approach has to be used wherever possible. It is self-evident that the shore-parallel storm waves hitting the coast involve higher hazard levels than shore-normal wave fronts. The wave breaker type is different in each sector: plunging waves are typical of the sector I and II, while collapsing and surging waves characterise the sector V and, the sectors III and IV respectively; the plunging waves break with more energy than the others. With regards to the exposure, the worst condition happens in the sectors I and II because they are exposed to the storms approaching from the N, NE and NW. The partial rating for each category of parameters taken into account and the total ratings for each coastal sector (CISA rating) are provided in Table III.

V. CONCLUSION

Coastal cliff retreat is difficult to assess and model due to the episodic nature of failures and the complexity of retreat mechanisms controlled by a number of factors dependent on the properties of rock masses and meteo-marine conditions. An ideal method for assessing cliff instability susceptibility along a coastal stretch needs for a preliminary calibration based on visual estimate, field inspections and cliff failure inventory. At same time, this method should be implemented with probabilistic and deterministic approaches. Nevertheless, it will present limitations and disadvantages in its application due to non-objective evaluations of the relative weight of the selected conditioning factors of future mass movements, difficulties of obtaining representative geotechnical data of rock masses, especially in karst areas, and because linked to a well-defined geological and environmental context.

The predictive capacity of the CISA method are not yet tested and, however, the results obtained in this study suggest that the procedure used may have a good potential for the assessment of the susceptibility of cliff failures in a typical Mediterranean coastal carbonate environment.
ACKNOWLEDGMENT

Many thanks to Riccardo Di Leva for the drawing in Fig. 7.

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Sex-specific effect of the LDL-receptor rs6511720 polymorphism on plasma cholesterol levels.
Results from the Czech post-MONICA study.

Jaroslav A. Hubacek, Vera Lanska, Vera Adamkova

Abstract—Aims: Acute coronary syndrome is one of the most common causes of death in industrialized countries. From the conventional risk factors, high plasma cholesterol and dyslipidemia are probably the most commonly analyzed. Plasma cholesterol is influenced equally both by environmental (physical activity, dietary habits) and genetic factors.

Methods: Rs6511720 (G → T) variant was analyzed by PCR-RFLP in 2,559 adult individuals (aged 28-67 years, Czech post MONICA study) with known plasma lipid parameters. ANOVA was used for statistical analyses.

Results: Genotype distribution was similar to the neighboring countries and genotype frequencies (0.085 for minor T allele) were within the Hardy-Weinberg equilibrium (P = 0.44). Polymorphism has no effect on plasma lipid levels in males. In females, carrier’s of the minor T allele has significantly lower concentration of total cholesterol (5.50 ± 1.07 mmol/L vs. 5.86 ± 1.16 mmol/L; P < 0.005) and LDL cholesterol (3.43 ± 1.00 mmol/L vs. 3.70 ± 1.05 mmol/L; P < 0.002). HDL cholesterol and triglycerides were not influenced by the analyzed polymorphism.

Conclusions: Rs6511720 (G → T) polymorphism within the LDL-receptor gene has an effect on plasma cholesterol levels in Czech females, but not in males.

Keywords— Czech population, LDL-receptor, plasma lipids, polymorphism.

I. INTRODUCTION

Coronary artery disease (CAD) and particularly acute coronary syndromes (ACS) are among the most frequent causes of death in developed industrial countries. Significant proportion (but definitely not all patients) of the ACS patients exhibit one on more from the five conventional risk factors (smoking, hypertension, type 2 diabetes, overweight/obesity and dyslipidemia) [1]. Dyslipidemia is characterized by high plasma total cholesterol and/or low density lipoprotein (LDL) - cholesterol, high levels of plasma triacylglyceroles (TG) and low levels of plasma high density lipoprotein (HDL) – cholesterol. Especially elevated plasma LDL cholesterol levels are associated with the ACS incidence [2].

High levels of plasma cholesterol result both from the unfavorable environmental factors (low physical activity and unhealthy dietary habits – low intake of fruits and vegetables and high intake of saturated fats) and from genetic predispositions (both in form of the common polymorphisms and rare mutations). It is estimated that both environmental and genetic factors influence final plasma cholesterol levels from 50%.

Seek for the genes (and variants) associated with plasma cholesterol levels was very successful in the case of genome wide association studies (GWAs). Couple of almost simultaneously published papers introduced dozens of variants associated with plasma lipid parameters [3]-[5]. Despite the fact, that the results were replicated on independent population, there is one major weakness of such studies. They have been performed mainly on west European populations with similar environmental influence and further replications are necessary [6].

Among the variants detected through the GWAs approach, are also variants within the LDL – receptor gene. LDL – receptor (OMIM acc. No. 606945) is the cell surface receptor which plays a very important role in cholesterol homeostasis. Binding to the apolipoprotein B, LDL – receptor internalize the cholesterol rich LDL – particles. One of the variant within the LDL – receptor gene, which was suggested by GWAs to influence significantly plasma cholesterol levels is the intronic rs6511720 polymorphism (G → T exchange).

To replicate the original finding, we have analyzed, if rs6511720 polymorphism at LDL – receptor locus is associated with plasma lipids in the adult Czech Slavonic population.
II. MATERIAL AND METHODS

A. Analyzed individuals

In our study, 2,559 adult (29 - 68 years at the time of examination at 2000/2001, mean age 52.0 ± 10.2 years) individuals selected according to the WHO MONICA Project protocol were examined and analyzed [7]. All participants were of Caucasian ethnicity and all signed the informed consent. 1,191 males and 1,368 females were included in the study. The study was approved by the institutional Ethics committee and conducted according to the Good Clinical Practice guidelines.

B. Laboratory and clinical analyses

DNA was extracted from peripheral blood white cells. Rs6511720 was genotyped (for more details see [8]) using the nested polymerase chain reaction and restriction fragment length polymorphism (PCR – RFLP) analysis using the oligonucleotides 5’aac atc aca ttc tca gcc atc ceg g and 5’ ttg tag aag tga gtt ctc get tgg. Product from the first PCR was diluted 1 : 50 and used as a template for the second amplification with oligonucleotides 5´ acc ggg gat gat gat gat and 5´  ttg cct aag act tcc tta aca ttt g. Restriction enzyme MboI was used to distinguish the alleles. All chemicals were purchased from Fermentas (Burlington, Ontario, Canada). Common G allele was represented by uncut product (132 bp) while the presence of restriction fragments of 107 and 25 bp, represented the minor T allele.

C. Biochemical analyses

Lipoprotein parameters (assessed in plasma after an overnight fast) were measured using standard enzymatic methods in CDC Atlanta accredited laboratory.

D. Statistical analyses

The Hardy-Weinberg test (http://www.tufts.edu/~mcourt01/Documents/Court%20lab%20-%20HW%20calculator.xls) was applied to confirm the independent segregation of the alleles. ANOVA was used for 0-% HW calculator.xls) was applied to confirm the significant association between the genotypes and plasma lipids. For LDL – cholesterol, there was a trend to the lower concentrations in T allele carriers.

Among males (Table II), we have not detected any significant association between the rs6511720 polymorphism and plasma lipids. For LDL – cholesterol, there was a trend to the lower concentrations in T allele carriers.

In contrast, variant was associated with plasma cholesterol levels in females. Both plasma concentration of total cholesterol (P < 0.005) and plasma LDL – cholesterol (P < 0.002) were significantly higher in common GG homozygotes than in carriers of the minor T allele (for more details, see Table III). Plasma levels of HDL – cholesterol and triacylglycerols were not influenced by the rs6511720 polymorphism.

The significance of the obtained results has not changed (neither in males, nor in females) if adjustment for body mass index, smoking status and age was performed. The finding, that the rs6511720 polymorphism within the LDL – receptor gene influence plasma lipids in females but not in males is not in agreements with originally published results.
[3]-[5], however, we have already confirmed some [9] but not all [10] [11], results from GWAs studies. Surprisingly it is a common feature of genetic association studies. Positive results published in high impact journals are more likely to have high bias scores (and smaller sample sizes) [6] [12] and thus are more prone to present incomplete and also false positive results.

Our results underline the importance of the replication studies, especially if they are performed on some clearly distinguished population subgroups (males vs. females, smokers vs. never smokers, etc.). They can point on some context dependent effect of the genetic polymorphisms on analyzed biochemical or anthropometrical parameter. Also the potentially different genetic background or different environment could play an important role in expression of the effects of the different alleles.

ACKNOWLEDGMENT

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Segmentation of Brain MRI Image Based on Clustering Algorithm

Siti Noraini Sulaiman, Noreliani Awang Non, Iza Sazanita Isa, Norhazimi Hamzah

Abstract— Medical images are widely used by the physicians to find abnormalities in human bodies. The physicians used the findings to plan further treatment for the patient. However the images sometimes are corrupted with a noise which normally exist or occurs during storage, or while transfer the image and sometimes while handling the devices. Therefore the need to enhance the image is crucial in order to improve the image quality. Segmentation technique for Magnetic Resonance Imaging (MRI) of the brain is one of the method used by radiographer to detect any abnormality happened specifically for brain. The method is used to identify important regions in brain such as white matter (WM), grey matter (GM) and cerebrospinal fluid spaces (CSF). In this project, the image segmentation via clustering method is used to cluster or segment the images into three different regions which represent the white matter (WM), grey matter (GM) and cerebrospinal fluid spaces (CSF) respectively. These regions are significant for physician or radiographer to analyse and diagnose the disease. The clustering method known as Adaptive Fuzzy K-means (AFKM) is proposed to be used in this project as a tool to classify the three regions. The results are then compared with fuzzy C-means clustering. The segmented image is analysed both qualitative and quantitative. The results demonstrate that the proposed method is suitable to be used as segmentation tools for MRI brain images using image segmentation.

Keywords— Image Processing, Image Segmentation, Brain MRI image, Clustering.

I. INTRODUCTION

Medical image normally used by the physicians to detect abnormalities in body system. It is also used for the treatment planning. Various medical images techniques used to sense the irregularities in human bodies such as Magnetic Resonance Imaging (MRI), Computerized tomography (CT), and Ultrasound (US) imaging. In such a case, the radiographer used a tool to make the decision of the medical images analysis easier. It also helps radiographer make accurate decision about the corresponding image. The radiologist use medical image to identify the tumours, tissues, and its anatomical structures [1]. But there are many problems faced when performing MRI procedure. The problems are the image generally have non-linear characteristics and sometimes are corrupted with noise. These problems make the radiologist faced difficulties in identifying of tumors, tissues and its location as well as difficulties to study the anatomical abnormal growth of glands. Finally, these may lead to inconveniences in making decision [1].

Many segmentation methods have been introduced in the literature [2]. In digital image processing, segmentation refers to the process of splitting observe image data to a serial of non-overlapping important homogeneous region [3]. Clustering algorithm is one of the process in segmentation. In the analysis of medical images for computer-aided diagnosis and therapy, a preliminary processing task often required is segmentation [3, 4]. Besides that, by using computer aided, image processing can be applying image reconstruction. It is very important to medical field because radiologist can identify the abnormality happen at the brain. Since radiologist can determine abnormality in the patient brain, they can give the best treatment for the patient.

There various image segmentation techniques based on clustering. For examples of clustering algorithm are K-means (KM) clustering, Moving K-means (MKM) clustering and Fuzzy C-means (FCM) clustering. Clustering is the process of separating data into group of similarity [5]. It also known as procedure of organizing objects into groups whose members are similar in certain way, whose goal is to identify structures or clusters existing in a group of unlabelled data[6]. Clustering algorithm are normally being used in computer, engineering and mathematics field [7]. In the past few decades, the uses of clustering algorithm have been broadening to medical fields, due to the development and advancement of medical imaging fields. Examples of medical images are image of brain, bone, and also chest. Clustering algorithm is suitable in biomedical because it will make the analysis easier.

Segmentation via clustering can also be used to detect the three regions at the brain image. Magnetic Resonance Image (MRI) of brain is one of medical imaging tools used to detect abnormality in brain. From the MRI brain images, the radiologist normally interested to look for three significant regions. The three regions are white matter (WM), grey matter (GM) and cerebrospinal fluid spaces (CSF) [3, 6]. Figure 1 shown three regions of normal MRI brain image. The precise
measurement of these three regions is important for quantitative pathological analyses and so becomes a goal of lots of method for segmenting MRI brain image data.

![Normal Brain MRI Image](image)

**Figure 1:** The normal brain MRI image

In this paper segmentation via clustering method named Adaptive Fuzzy K-means (AFKM) clustering is used to segment the MRI brain image into three different regions. The AFKM method is proposed to prove that it can classify and segment the MRI brain image better than conventional method. AFKM clustering algorithm is combination of KM, MKM and FCM clustering. The features of AFKM are to provide a better and more adaptive clustering process.

The remaining of the paper is organized as section II presents about AFKM clustering algorithm is proposed. Section III describe data analysis of MRI brain image. Section IV presents a result analysis comparison of performance of FCM and AFKM clustering algorithm. Besides that, the comparison detail of qualitative and quantitative also presented. Lastly, Section V concludes of this paper.

## II. METHODOLOGY

In medical field, Medical Resonance Image (MRI) is one of the methods used to detect abnormalities in human body. The clustering algorithm for image segmentation was introduced to the MRI brain image in order to segment the image.

In this paper, a new method of clustering algorithm based segmentation known as technique is recommended [6]. The segmentation technique used to be implemented medical image like MRI. It use to exquisite soft tissue contrast between normal tissue and pathologic tissue. The proposed method for this paper is then comparing with conventional method known as Fuzzy C-means (FCM). In this section algorithm FCM and AFKM are briefly discussed in II(A) and II(B).

### A. Fuzzy C-means (FCM) clustering Algorithm

Bezdek is the person who introduced Fuzzy C-means (FCM) algorithm. The FCM is one of the most commonly used clustering algorithm [6, 7]. FCM clustering is constructed based on the same idea of definition cluster centers by iteratively regulating their locations and minimizing an objective function as K-Means (KM) algorithm [6]. The advantage of FCM is, it allows more flexibility when dealing with multiple cluster by introducing multiple fuzzy membership grades [6].

### B. Adaptive Fuzzy K-Means (AFKM) clustering Algorithm

In this paper, AFKM method is recommended to be used to process MRI images. It is the latest type of clustering algorithm proposed by [7]. The AFKM is combination of fundamental theories of conventional K-means and MKM clustering algorithm (i.e., assigning each data to its closet centre or cluster) and the conventional Fuzzy C-means (FCM) clustering algorithm (i.e., allows the data to belong to two or more clusters or centres). The objective function of AFKM is calculated using the equation:

$$J = \sum_{k=1}^{n} \sum_{t}^{N} \left( \frac{M_{kt}^m}{v_t} \right) \|v_t - c_k\|^2$$

where $M_{kt}^m$ the fuzzy membership function and $m$ is the fuzziness exponent. The degree of being in a certain cluster is related to the inverse of the distance to the cluster. The new position for each centroid is calculated using the equation:

$$C_k = \frac{\sum_{t=1}^{N} (M_{kt}^m) v_t}{\sum_{t=1}^{N} (M_{kt}^m)}$$

where,

$$(M_{kt}^m) = M_{kt}^m + \Delta M_{kt}^m$$

where $(M_{kt}^m)$ is the new membership and is defined as:

$$\Delta M_{kt}^m = \alpha(c_k)(e_k)$$

and $e_k$ is error of belongingness. Then, the value of $e_k$ is calculated by

$$e_k = B_k - \hat{B}_k$$

The AFKM algorithm improved the clustering with the introduction of belongingness concept where it measures the degree relationship between centre and its members. The degree of belongingness, $B_k$ is calculated using:

$$B_k = \frac{e_k}{M_{kt}^m}$$

The objective is to minimize the objective function from equation (1). The process is repeated iteratively until the center is no longer moved all data have been considered.
Images of brain MRI are obtained from internet database. The images are processed with AFKM and FCM clustering algorithm and comparison is made between the two clustering algorithms. The flow chart for the whole process is depicted in Figure 2.

III. DATA ACQUISITION AND ANALYSIS

The method Adaptive Fuzzy K-means (AFKM) clustering algorithm is introduced to segment a MRI brain image but usually the MRI brain image used computer-aided to detect any irregularities happened. In this paper, six of MRI brain images obtained from internet databases are chosen to be tested the AFKM algorithm, as shown in Figures 3(a) until 3(f) respectively.

To implement the performance analysis, qualitative and quantitative are considered. There are three evaluation functions used in quantitative analysis obtain from Liu and Yang [8]:

\[
F(I) = \sqrt{R \sum_{i=1}^{R} \frac{e_i^2}{A_i}} \tag{7}
\]

Equation \(F'(I)\) and \(Q(I)\) are proposed by Borsotti [9] where the \(Q(I)\) is the improved version of \(F(I)\). The equations are gives follow:

\[
F'(I) = \frac{1}{1000(NxM)} \sqrt{\sum_{A=1}^{Max} [R(A)]^{\frac{1}{R}} \sum_{i=1}^{R} \frac{e_i^2}{A_i}} \tag{8}
\]

\[
Q(I) = \frac{1}{1000(NxM)} \sqrt{R \sum_{i=1}^{R} \left[ \frac{e_i^2}{1 + \log A_i} + \left( \frac{R(A)}{A_i} \right)^2 \right]} \tag{9}
\]

For evaluation of the cluster quality, the most fundamental benchmark is the mean squared error (MSE). It could be described as follows:
IV. RESULT AND DISCUSSION

From these images, the performance analysis of qualitative and quantitative are implemented. The qualitative analysis depends on the human visual. Human visual can interpret the images based on capability and segmentation algorithm of conventional method like FCM and the new method proposed which is AFKM. It can detect the region of interest like GM, WM and CSF. For quantitative analysis, it refers to the performance of segmentation of the image. It produces by proposed algorithm. The conventional algorithm will compared with a new proposed algorithm. The result of quantitative analysis taken based on three evaluation functions. The three functions of quantitative analysis are F(I), F'(I) and Q(I). The image size can calculate from N x M. For evaluation of the cluster, the mean squared error (MSE) is the one most fundamental benchmark. Besides that, these functions related more to the visual judgment. For the better result of segmentation, AFKM values of F(I), F'(I) and Q(I) are smaller than FCM values. Both of result of qualitative and quantitative will be presented in section IV (A) and IV (B).

A. Qualitative Analysis

For the result in qualitative analysis, six images are used. Qualitative analysis is to examine usually whether the resultant image is good or not. The performance is examine visually in qualitative analysis. The segmentation performances are compared with conventional methods of FCM and new method proposed of AFKM. Clustering algorithm used in this paper is to segment the MRI brain image into here regions i.e. the GM, WM and CSF, therefore the clustering algorithm is chosen to have three clusters. The result is then compared with FCM algorithm. From the result shown in Figures 4 to 5, it can be observed that quality of image is not perfect compared to the AFKM method. The weakness of FCM method is it over segment the image which leads to image become too bright. But using AFKM, it can segment the image clearly and the region of interest is sharper.

\[
MSE = \frac{1}{N} \sum_{j=1}^{M} \sum_{i \in s_j} \| x_i - c_j \|^2
\]  

(10)

Figure 4: The image 1 of segmentation image with three clusters: (a) Original image. (b) FCM. (c) AFKM

Figure 5: The image 2 of segmentation image with three clusters: (a) Original image. (b) FCM. (c) AFKM

Figure 6: The image 3 of segmentation image with three clusters: (a) Original image. (b) FCM. (c) AFKM

Figure 7: The image 4 of segmentation image with three clusters: (a) Original image. (b) FCM. (c) AFKM

Figure 8: The image 5 of segmentation image with three clusters: (a) Original image. (b) FCM. (c) AFKM

After implementing AFKM algorithm, the image looks clearly in the visual compared to the conventional method of FCM. The resultant images are shown in Figures. 6 and 7. By FCM, the MRI brain image is brighter compared to AFKM. It happens because the FCM have over segment of the image. It can give effect on segmentation and three regions cannot be
detected clearly. The images of AFKM become sharper and clearer.

Figure 9: The image 6 of segmentation image with three clusters: (a) Original image. (b) FCM. (c) AFKM

By using FCM, the image is unclear. It is because the image becomes bright and not meets criteria of segmentation. But when AFKM method applied, the images are sharp and the segment of WM, GM and CSF are correctly. The resultant images are shown in Figures. 8 and 9 respectively.

From the resultant images are shown in Figures. 4 to 9, a new method proposed of AFKM can give better performance of segmentation technique compare the conventional method of FCM.

B. Quantitative Analysis

The analysis of quantitative is evaluated based on the three benchmark functions. It is also mentioned in section III. The analysis also evaluates by fundamental benchmark is mean squared error (MSE). The quantitative analysis is to support the qualitative finding in section III(A). The result of quantitative analysis show in Tables 1 to 4. These tables summarize the segmentation of the quantitative estimation. All the result gets from a comparison of FCM and AFKM clustering method. From the comparison, new method of AFKM produces the better result compared to conventional method of FCM. The new method proposed of AFKM produce the smaller values of all MSE, F(I), F'(I) and Q(I) analysis. So, it can be conclude that the AFKM method is successful segmentation. It is because AFKM can detect the three regions at MRI brain. In addition, the proposed AFKM manages to segment the image successfully with less noisy pixel. Generally, these interpretations specify that the AFKM might be a better methodology in terms of image segmentation application.

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Note: For the rest of the Tables, bolded numbers show the best result obtained for each analysis.

V. CONCLUSION

In this paper, the new method of AFKM clustering algorithm is present. The AFKM clustering is combination with MKM, KM and FCM. The result can prove that using AFKM can get sharper and clearer of segmentation in MRI brain image. This process is a good method for segmentation. The technique was used to segment the three regions in MRI brain image using clustering algorithm. The results get from the qualitative and quantitative of MRI brain image. In the future, AFKM method can apply in engineering field, agriculture field and also nutrition field.

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BIographies

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Effects of adolescent idiopathic scoliosis on postural balance and muscle activity

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Abstract—In adolescent idiopathic scoliosis, it is very important to understand the effects of idiopathic scoliosis on postural balance and muscle function. In this paper, we assessed the correlation between inclination angle and muscle activity in static and dynamic sitting condition. Twelve subjects with idiopathic scoliosis participated in this study. Inclination angle was measured in sagittal and frontal plane by unstable board with accelerometer. Surface electrodes were attached to the external oblique, thoracic erector spinae, lumbar erector spinae, and lumbar multifidus muscles, bilaterally. In static sitting condition, subjects were instructed to sit the board comfortably for 30 seconds and then neutral tilting angle was measured. And, in dynamic sitting condition, tilting angle and muscle activity when pelvic anterior, posterior, left, and right pelvic tilting induced by the structure of the board was analyzed. The results show that mean inclination angle in posterior and left side more increased significantly than other side in static and dynamic sitting condition. There was a difference in muscle activity pattern between both sides of abdominal and erector spinae muscles. Erector spinae muscles of subjects with AIS more increased on right side than left side when pelvic anterior, posterior, and left tilting. Especially, thoracic erector spinae muscle on right side increased in all direction. From these results, we suggested that adolescent idiopathic scoliosis affect the postural asymmetry and muscular imbalance in sitting. Furthermore, we confirmed that unstable board can facilitate to evaluate the correlation efficiently between trunk asymmetry and abnormal muscle activity caused by idiopathic scoliosis.

Keywords—Adolescent idiopathic scoliosis, inclination angle, muscle activity, unstable board.

I. INTRODUCTION

A DOLESCENT idiopathic scoliosis (AIS), which is defined as abnormal lateral curvature of the spine, is the most common deformity occurring about 4% of children between the ages of 11 and 17 years of age [1]-[4]. AIS is found more frequently in females than males (sex ratio of 8:1) [5]. Although there is no clear evidence, idiopathic scoliosis has been associated with progressive loss of muscle function and balance in the sagittal and frontal plane in standing and sitting.


Although there is evidence that how idiopathic scoliosis affect the entire skeletal system in our daily life, the correlation between postural balance and muscle activity in sitting has not been studied. It is very important to understand the association between sitting posture and muscle activity caused by idiopathic scoliosis in static and dynamic sitting condition due to sitting has become the most common posture with the increase in study time of students.

Accordingly, we proposed the efficient method for evaluating the effect of idiopathic scoliosis on postural balance and muscle activity using unstable board. Also, we evaluated the
characteristics of inclination angle and EMG pattern of AIS patients for preventing degenerative scoliosis and providing useful information to treat idiopathic scoliosis in rehabilitation medicine.

II. METHODS & MATERIALS

A. Subjects

Twelve female subjects with adolescent idiopathic scoliosis were included in this study. All subjects were recruited from the Department of Rehabilitation of Chungnam National University Hospital in Daejeon, Republic of Korea. The mean age, height, body weight was 15.67±2.23 years, 16.03±4.58 mm, and 49.83±6.43 kg, respectively. The inclusion criteria for patients were posteroanterior full spine standing X-ray evidence of idiopathic scoliosis with a lumbar or thoraco-lumbar curve and no previous conservative or surgical treatment for the scoliosis. The mean cobb angle was 21.6° (range: 16° to 29°) and major curve defined by a Cobb angle was to the right (convex side) as shown in Fig. 1. Exclusion criteria were pervious orthopedic surgery, central or peripheral neurologic disorders, or other spinal disorders. All subjects were informed a full explanation regarding the protocol and provided written consent prior to their participation.

B. Experimental Setup

To assess the postural balance of subjects in sitting, we used unstable board with 3-axis accelerometer which shape and appearance was hemisphere. As shown in Fig. 2, curvature radius of the board was designed to incline in the sagittal (anterior-posterior, AP) and frontal (left-right, LR) plane for evaluating the postural balance of trunk and pelvis as well as muscle activity according to trunk and pelvis movement. Accelerometer which was positioned to middle bottom of the board facilitated measurement on asymmetry sitting posture [15].

Muscle activation patterns were recorded using the Noraxon Telemyo 2400T (Noraxon Inc., Scottsdale, USA). Wireless surface electrodes (Noraxon Inc., Scottsdale, USA) were attached to the external oblique (just below the rib cage, along a line connecting the most inferior costal margin and the contralateral pubic tubercle, EO) [16], thoracic erector spinae (5-cm lateral to the T9 spinous process, TES), lumbar erector spinae (3-cm lateral to L3 spinous process, LES) [17], and lumbar multifidus (L5 level, parallel to a line connecting the posterior superior iliac spine and L1-L2 interspinous space, LM) muscle [18] bilaterally as shown in Fig. 3. Before placing electrode, the skin was shaved and cleaned with alcohol to reduce skin resistance.

C. Experimental protocol

Experimental protocol was divided into two conditions: static and dynamic sitting. First, in static sitting condition, subjects were instructed to sit in their usual manner on the board, which is located in the center of stool, with their arms crossed on contra-lateral shoulder for 30 seconds. And, lastly, in dynamic sitting condition, subjects were asked to perform pelvic anterior, posterior, left, and right tilt using curvature structure of the board, and then subjects keep the posture for 5 seconds, respectively.

From these experimental conditions, we evaluated the neutral sitting posture and correlation between postural balance and muscle activation caused by idiopathic scoliosis. A foot support was used to prevent the influence of leg movement, and it was adjusted to support the feet by keeping and ankle angles at 90° [19]. Before the experiment, all subjects had enough time to adapt to instability by inducing the board.

D. Data Analysis

Inclination angles with trunk and pelvic movement in frontal and sagittal plane under two conditions were analyzed, sampled at a rate of 100 Hz, using LabVIEW 2010 (National Instrument CO., Texas, USA).

All EMG signals from trunk muscles were amplified, bandpass filtered (passband 20-450 Hz), notch-filtered at 60 Hz, and then sampled at 1000 Hz. MyoResearch Master XP 1.07 (Noraxson Inc., scottsdale, USA) was used to analyze the measured data.
To normalize the difference of muscle contraction for individuals, EMG data was expressed as a percentage relative to the maximum voluntary contraction (MVC). Two MVC trails were conducted to obtain the maximal electromyography of abdominal and erector spinae muscles. In this procedure, subjects attempted to flex and/or extend the upper trunk in the sagittal plane with the maximum effort, and then hold it for five seconds whereas examiner pushes down shoulders of subject [20], [21].

Statistical analysis was performed using SPSS 18.0 software (SPSS Inc., Chicago, USA). Independent t-test was used to examine the difference in angle variation and pressure distribution between PA and PS group, at p < .05 level.

III. RESULTS

A. Inclination Angle

In static sitting condition, mean inclination angle in AP and LR direction was -1.58° and -1.65°, respectively, as shown in Fig. 4. Angle variation was tilted to posterior and left side.

In dynamic sitting condition, mean inclination angle between anterior and posterior as well as left and right was compared to each other as shown in Fig. 5. Anterior and posterior tilting angle was 5.46° and 5.50°, respectively. In addition, left and right tilting angle was 6.73° and 3.46°, respectively. Inclination angle in posterior and left side was larger than other direction, and inclination angle in left side increased more significantly compared to that of side (p < 0.001).

B. Muscle Activity

Differences in muscle activity of both sides during pelvic anterior tilting are presented in Fig. 6. External oblique, thoracic erector spinae, lumbar erector spinae muscle increased on right side whereas lumbar multifidus was reduced on that side. In contrast, all trunk muscles more increased on right than left side when pelvic posterior tilting as shown in Fig 7. Especially, differences in muscle activity of thoracic erector spinae, lumbar erector spinae, and lumbar multifidus increased significantly on right side during pelvic left tilting (p < 0.001, p < 0.000, and p < 0.001, respectively) as shown in Fig. 8. Muscle activity of external oblique, lumbar erector spinae, and lumbar multifidus increased on left side whereas thoracic erector spine only decreased on that side when pelvic right tilting as shown in Fig. 9.
balance and muscle activity in sitting. Inclination angle was tilted to posterior and left side in static and dynamic sitting. And, differences in mean inclination angle between pelvic left and right tilting was more increased significantly than anterior and posterior tilting. AIS patients showed larger muscle activity on the convex side of thoracic erector spinae, lumbar erector spinae, and lumbar multifidus. And, these muscular imbalance of erector spinae muscles are associated with inclination angle which was tilted to posterior and left side. Consequently, we concluded idiopathic scoliosis cause postural asymmetry and muscular imbalance in sitting. This paper suggested that structure of unstable board may be utilized to assess the effect of postural and muscular imbalance caused by low back pain, scoliosis, leg length discrepancy, pelvic asymmetry on the individual’s activities in daily life. Future research is required to compare the inclination angle and muscle activity between AIS patients and control subjects.

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The Double Reflection Control of Direct Solar Light in Deep Atrium Type Spaces

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Abstract—Redirecting sunlight by double reflection in deep atrium type space allows a better control of natural lighting to adjacent spaces. The problem resulting from a permanent change in sun angle to the horizontal plane and their azimuth. A possible solution is to use pairs of flat horizontal rotary shafts contained in a cylinder that rotates after the sun azimuth. This way we can get in the atrium the sunlight always vertical, almost vertical or scattered as desired. Additionally, by positioning the panels at specific angles can be obtained on the winter nights insulation and natural ventilation and protection from the sun in hot summer days.

Similarly, on days without sunshine is allowed to spread between diffuse light and in temperate nights the space can be ventilate.

Keywords—atrium, deep space, direct sunlight reorientation, double reflection, renewable energy.

I. INTRODUCTION

Using the atrium in a building concept could increase its efficiency by using the renewable energy, thus reducing its carbon footprint.

The atrium creates an identity, a more interesting space and it stimulates the interpersonal relations. In addition to this, the atrium can use the natural light, can ventilate the building and can contribute to the heating or cooling of the air.

Depending on the conformation, there are some improvements to be brought in order to increase the level of comfort and to decrease the level of used energy.

In high buildings, the height and width ratio of the atrium increases, thus rising to lighting problems towards the lower levels, progressively.

Inserting the natural light vertical, a great part of the artificial lighting can be replaced in a controlled way, the latter being a big energy consumer.

II. PROBLEM FORMULATION

The constant changing of azimuth and angle in vertical plane of direct solar rays creates glare problems, unequal lighting, local superheating at higher levels and poor illumination at lower levels (Fig 1).

As seen below, in The Office Building in Barcelona [1] (Fig. 2), there is a risk of accidental direct sunlight strips to enter in the offices. In Stuttgart’s Art Museum (2004, arch. Hascher Jehle) [2] (Fig. 3), rotating panels are used to control the sunlight. In the Youth Forum (YUFO) [3] (Fig.4) in Möglingen (Germany, arch. P. Hübner 1993), the “Solar Eye” follows the sun in winter and turns 180° away in summer.
A solution is to dissociate the apparent movement of the sun in the sky in two components:

1. The horizontal component
2. The vertical component

1. The horizontal component could be followed by using horizontal spinning panels with a mirror surface in a cylinder that rotates by the sun so as to keep the panel's axis always perpendicular on the direct sun rays (Fig. 4, Fig. 5 and Fig. 6).

2. The vertical component could be followed by rotating the panels around the horizontal axis.

Due to the geometry of the system with a single row of panel, its efficiency is low (shown by vertical stripes "V" Fig7); a great deal of the sun rays strips (marked by "G" from glare) cannot be redirected (and Fig. 8).

For panel of 100 cm. width and 60° angle of incident sun rays result 25.88 width vertical rays (29.88% from incident rays) and 54.93 width "glare rays" (63.43% from incident rays) (Fig. 9)
The proposed solution is the double redirecting (or the double reflection), by mounting two panels on an axis, one is attached to the axis (the primary panel) and the other one is mobile towards the axis (the secondary panel).

Next are shown the panels positions for three particular angles (60°, 45° and 30°), with the efficiency measured geometrically (Fig. 10, Fig. 11 and Fig. 12). Under 30° the secondary panel is out of service (Fig. 13).

Fig. 10 System under 60° incident sun rays

Fig. 11 System under 45° incident sun rays

Fig. 12 System under 30° incident sun rays

Fig. 13 System under 20° incident sun rays

Table 1

<table>
<thead>
<tr>
<th>Incident rays</th>
<th>dimension (cm)</th>
<th>%</th>
</tr>
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<tbody>
<tr>
<td>Incident rays 60°</td>
<td>86.6</td>
<td>100</td>
</tr>
<tr>
<td>Reflected rays</td>
<td>52.03</td>
<td>60.08</td>
</tr>
<tr>
<td>Losses</td>
<td>34.57</td>
<td>39.92</td>
</tr>
<tr>
<td>Efficiency under skylight</td>
<td>34.57/100</td>
<td>34.57</td>
</tr>
<tr>
<td>Incident rays 45°</td>
<td>70.71</td>
<td>100</td>
</tr>
<tr>
<td>Reflected rays</td>
<td>47.35</td>
<td>66.96</td>
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<tr>
<td>Losses</td>
<td>23.36</td>
<td>33.04</td>
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<tr>
<td>Efficiency under skylight</td>
<td>47.35/100</td>
<td>47.35</td>
</tr>
<tr>
<td>Incident rays 30°</td>
<td>50</td>
<td>100</td>
</tr>
<tr>
<td>Reflected rays</td>
<td>36.3</td>
<td>72.6</td>
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<tr>
<td>Losses</td>
<td>13.7</td>
<td>27.4</td>
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<tr>
<td>Efficiency under skylight</td>
<td>36.3/100</td>
<td>36.3</td>
</tr>
<tr>
<td>Incident rays 20°</td>
<td>34.17</td>
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<td>Reflected rays</td>
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<tr>
<td>Losses</td>
<td>6.85</td>
<td>20.05</td>
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<tr>
<td>Efficiency under skylight</td>
<td>27.32/100</td>
<td>27.32</td>
</tr>
</tbody>
</table>

A summarization of the data:

It can see that the system perform better around 45° because the real efficiency is the one under the skylight.
The sun position in the sky referred to a fixed point on Earth is known on any time, on any date in the past, present or future. Therefore the adjustment of the three subsystems: cylinder, primary panels and secondary panels becomes a space geometry matter depending on time.

The computer does the adjustments by a dedicated software having the precise location input; but there is a possibility to choose other pre-established scenarios depending on some activities.

The position update is made at regular periods of time (once every 1... 10 minutes) depending on the required accuracy and the building geometry.

The system is activated by actuators powered by solar panels on the roof, one is for the cylinder, one is for the primary panels (plus the gearing system for all the primary panels) and one is used for each secondary panel.

In addition, the cylinder has ventilating hatches, electrically controlled, in order to insure the ventilation depending on the main direction of the wind.

The primary panels include thermal insulation; in a horizontal position, they isolate additionally the atrium during cold nights.

The dynamic mapping of the clouds in the sky (part of an on-going study) would co-operate interactively with the computer which directs the panels to get the maximum of the natural light (the panels are complete open on cloudy days and they are dynamically closed and open on partially sunny days). A short-term prediction of the clouds position is required to have the smallest fluctuations in the indoor lighting. Of course, the fluctuations are lowered by a dimmable lighting system.

There are different scenarios witch can be applied in different situations such as: "hot summer day" when a major part of direct sun light can be reflected outside by mirror surface leaving only the diffuse light inside, "summer night" when the hatches are opened for ventilation (Fig. 14), "winter night" when the panels are in horizontal position in order to insulate the atrium (Fig. 15), "cloudy spring day" when the hatches opened for ventilation and the panels are in 90° position for diffuse light, Instead of having no panels (Fig. 16), "normal" (Fig. 17), "playful" (can be also dynamic) (Fig. 18), and so on.

IV. CONCLUSION

Controlling more accurate the movement of light in deep spaces allows substantial contribution to reducing or even eliminating artificial lighting at cost price given only by the acquisition of the system and installation of photovoltaic panels that power them. The remaining photovoltaic panels form an "electric farm" its production can be used internally or injected into the electric network.

V. CLAIMS

The double panel system with different orientation though related of each panel.
Operating scenarios:
- day / night,
- winter /spring-autumn / summer,
- various occasional activities.
Panel structure with insulation used for atrium's thermal insulation in cold nights.
The geometric "measure" of efficiency.

REFERENCES

Possible Assessment on Sustainability of Slopes by Using Electrical Resistivity: Comparison of Field and Laboratory Results

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Abstract—The stability of both natural and engineered slopes has always been a major concern of geotechnical engineers especially in recent years where an increasing number of development are encroaching into hillside areas. In tropical regions, soil slope failures due to frequent rainfall are quite common. Malaysia for example has witnessed various slope failures/landslides causing extensive loss in properties and even to the extent of causing casualties in the last two decades. These social and economic losses due to slope failures could be well minimized if identification of hazards/risks were established in early stages. One of the long term objectives of this whole research is to implement a quick method of assessing the factor of safety (FOS) in slopes by replacing the conventional soil parameters such as cohesion and internal angle of friction with electrical parameters such as resistivity. However, this paper is limited to the preliminary comparison of laboratory results obtained from controlled laboratory soil samples and results obtained from actual field samples. Results from both the laboratory controlled samples and actual field samples shows consistencies in the correlation between friction angle and electrical resistivity while correlations between moisture content and electrical resistivity shows a similar trend of decreasing moisture content with increase of electrical resistivity value.

Keywords—electrical resistivity, correlation, shear strength, slope stability.

I. INTRODUCTION

For the past two decades or so, Malaysia has witnessed many failures in slopes and landslides causing extensive loss in properties and even to the extent of causing casualties. With the current implementation of risk management system and standard operating procedure (SOP) for sustainable hillside development developed by the respective parties, among the required elements to be implemented in the afore mentioned systems are:

i. the identification of danger

ii. quantifying hazard

iii. determination of risk

One of the essential aspect to identify risk in slopes is to determine/calculate the factor of safety (FOS) which will indicate the stability of a certain slope. In the process of obtaining the FOS, among the crucial soil parameters to be obtained before calculating FOS are cohesion \((c)\), internal frictional angle \(\phi\), unit weight of soil \(\gamma\) etc. Since most of slope failures in Malaysia are mainly due to infiltration [1], the moisture content/pore water pressure also contributes to the FOS value. All these parameters are obtained for example through bore hole sampling.

In general practice, soil investigation (SI) incorporating bore hole sampling perhaps will produce the most reliable value of the relevant soil parameters for the purpose of actual calculation on factor of safety in slopes. However, bore hole sampling is in general time consuming and very expensive. Conventional methods of soil analysis mostly require disturbing soil, removing soil samples and analyzing them in laboratory where else electrical geophysical methods on the contrary allow rapid measurement of soil electrical properties such as electrical resistivity and conductivity directly from soil surface to any depth without soil disturbance [2].

Gue [3] mentioned that among the critical element in SOP is the need of checking of slopes especially in hillside development, which could be done among others by checking the FOS. For a regular checking and calculation of FOS in a certain stretch of slopes for the purpose of identification of risk/danger for example, bore hole sampling would not be practical due to the above mentioned reasons. This is because many bore holes are required to check the factor of safety at different locations on the slopes in order to determine the risk/hazard. Hence an alternate quick and less expensive method of assessing FOS is essential so as to enable rapid and extensive measurements and calculation of FOS at different points in slopes.

Therefore, this future quick assessment method which is based on electrical resistivity method is to preliminary check the factor of safety (FOS) of any slopes on initial and regular basis. Any slope could be checked and if the FOS falls within a certain range of a “prescribed values” which indicates high risk, a further confirmation of the FOS will then be conducted through the actual soil boring sampling or any other extensive method.
The general approach behind this quick assessment system is to eliminate the usage of physical soil parameters such as cohesion (c), internal frictional angle (φ) and unit weight (γ) as is currently being practice for the calculation of FOS and replace these physical parameters with their correlated electrical parameters such as resistivity. Therefore the simplified method at site will require a few steel rods implanted in the soil/slope serving as the electrodes, a reel of electrical wires and an existing multi meter to generate the factor of safety calculated through a set of empirical formula, charts and graphs all to be determined from this research.

The work of researchers in past and recent years have included correlation of electrical resistivity with various soil properties. Hassanein [4] have studied the relationship of electrical resistivity in compacted clay with hydraulic conductivity and some index properties. Earlier research had suggested the possible correlation of electrical resistivity with hydraulic conductivity which serves as a nondestructive mean for evaluating the quality of compacted soil liner [5]. An extensive work by Pozdnyakov [6] have looked into the effects of electrical resistivity in different soil types with varying water content, humus content, salt contents and several other parameters. Other researchers have also studied on estimation of water content of soil using electrical resistivity [7]. Others have used the knowledge of electrical resistivity to estimate liquefaction of soil [8], detecting and locating geomembrane failures [9], estimation of soil salinity for agricultural activities [10], etc. Syed et. al. [11] investigated the relationship of the electrical resistivity with soil parameters on homogeneous samples of sand, silt and clay at laboratory scale. Moisture content found to have strong relationship with resistivity. Poor correlations were observed between cohesion and friction angle with electrical resistivity for sand and silt samples, whereas clays samples showed a good correlations between shear strength parameters and resistivity. Syed and Zuhar [12] conducted some preliminary field work on actual slopes in the right direction.

In this paper a relook into the results obtained from previous laboratory and field works are conducted, compared and presented in order to highlight the feasibilities and uncertainties of this research. Recommendations are then proposed in order to keep the momentum going hopefully in the right direction.

II. MATERIALS AND METHODS

The research methodology consist of both field and laboratory investigations. The study area is located at University Technology PETRONAS, Perak, Malaysia as shown in Fig. 1. Field investigations comprise of electrical resistivity survey (VES) and soil boring. Laboratory investigations consist of soil characterization tests and electrical resistivity test.

A. Vertical Electrical Sounding

The vertical electrical sounding or 1D survey was conducted at the locations of boreholes (BH-01 to BH-10), using simple equipments and accessories in acquiring the electrical resistivity value e.g. handheld multimeter, D.C. power source, insulated wires, measuring tapes, stainless steel electrodes. The electrical sounding was conducted using Wenner electrode configuration with electrode spacing ranging from 0.5 to 3 meters. The apparent electrical resistivity of soil (ρa) is determined by equation (1).

\[ \rho_a = \frac{2\pi RL}{I} \]

The obtained apparent electrical resistivity values were inverted to true resistivity values using Ipi2win software and were used for interpretation. IPI2win is an open-source algorithm freely distributed by Moscow State University. The procedure for inversion involves automatic and manual technique. Initially automatic inversion were selected in order to get initial model and later on inversion models were refined or fine-tune using manual method until least RMS error was obtained.

B. Soil Boring

Soil boring was performed using percussion drilling set CobraTT equipped with 1meter core sampler. Depth of all boreholes (BH-01 to BH-10) was 3 meters. Prior to drilling PVC pipe was fixed in core sampler for easy and smooth recovery of soil samples from the core barrel. The obtained samples were then brought to the laboratory for soil characterization and electrical resistivity test in laboratory conditions.

C. Laboratory Tests on Index and Engineering Properties.

The basic idea behind this research is to estimate various soil properties using resistivity values. Therefore various soil characterization tests were performed to determine engineering properties of soil. Laboratory tests were performed on the soil samples obtained from boreholes, such as moisture content, unit weight, direct shear, sieve analysis, hydrometer test, liquid limit, plastic limit etc. as per methods suggested in British standards (BS).

D. Laboratory Resistivity Test on Field Samples

Electrical resistivity of soil samples from various depths was measured in order to determined resistivity values in
laboratory condition. Two disc electrodes were connected to both ends of cylindrical soil samples and also attached to DC power source and multi meter for current measurement. Potential difference varying from 30V, 60V, and 90V were applied and resulting variation in current were recorded. Laboratory temperature during electrical resistivity test was recorded as 24°C. Fig. 2 shows experimental setup for resistivity measurement in laboratory conditions.

![Fig. 2 Setup for laboratory electrical resistivity measurement](image)

The electrical resistivity of soil samples were determined by equation (2) and (3). Where \(V\) is voltage in volts, \(I\) is current in amperes, \(R\) is the resistance in ohms, \(A\) is the cross-sectional area of soil sample in meters, \(L\) is the length of soil sample in meters and \(\rho\) is the resistivity in ohms meter

\[
R = \frac{V}{I} \quad (2)
\]

\[
\rho = \frac{A}{L} R \quad (3)
\]

E. Laboratory Resistivity Test on Controlled Laboratory Samples.

Dry soils of separate sandy and silt size particles obtained from soil supplier were each mixed with 25%, 30%, 35% and 40% of distilled water in the laboratory. Mixing was done by means of a soil mixer and the samples were then left aside for at least 24 hours in the mixing bowl wrapped with plastic. Prior to the compaction process, the internal perimeter of the mold was lined with a thick plastic material. The specimens were then compacted in three equal layers using standard proctor hammer that delivers blows ranging from 15 to 45 blows per layer. The measurement of the electrical resistivity was done in the same manner conducted on the field samples.

III. RESULTS AND DISCUSSIONS

A. Soil Samples Description

A total of 79 soil samples were obtained from 10 boreholes (BH-01 to BH-10). Resistivity results and soil boring indicates different geological profiles ranging from silty sandy soils to sandy soils. Grain size analysis shows that soil samples from BH-01 to BH-06 are classified as “silty-sand” whereas soil samples obtained from BH-07 to BH-10 are mostly “sandy” soil samples according to British Soil Classification System (BSCS). Based on grain size distribution analysis, it can be concluded that 43 soil samples are silty-sand and 36 soil samples are course-grained sandy soils.

The controlled samples which were prepared in the laboratory as mentioned earlier were comprised of separate silt and sandy size particles and later data from the two separate sizes were combined for analysis.

B. Correlations Between Friction Angle and Electrical Resistivity

In order to establish relationship between electrical resistivity and various soil properties, simple regression analysis technique was used. Separate analysis for field and controlled laboratory samples was performed. The correlations between electrical resistivity and various properties of soil samples were evaluated using least-squares regression method. Linear, logarithmic, polynomial (quadratic and cubic), exponential and power curve fitting approximations were applied and the best approximation equation with highest correlation coefficient was selected.

The relationship between friction angle and field resistivity and laboratory resistivity (from samples obtained from field) indicates increasing logarithmic trend with \(R^2=0.33\) and \(R^2=0.25\) respectively for all soil samples as shown in Fig. 3 and Fig. 4. The obtained behavior is quite understandable due to the fact that shear strength of soil increases with decreasing moisture content [13]. This phenomenon could be clearly observed from the relationship between moisture content and electrical resistivity as shown in Fig. 5 and Fig. 6. The regression values of \(R^2=0.59\) and \(R^2=0.54\) respectively indicate the strong relationship between electrical resistivity and moisture content. Decrease in moisture content, increases electrical resistivity of soil. Hence, at higher resistivity values, higher friction angle will be observed. Another possible reason for this increasing trend might be due to the decreasing saturation of the soils. Since nearly saturated pores form bridges between particles and greater particle-to-particle contact [14], therefore decreasing saturation reduces particle-to-particle contact. Thus, lower and higher electrical resistivity associated with friction angles are results of decreasing or increasing electrical conductivity or resistivity in the pores and along the solid surface.

It should be noted here that besides being dependent to the amount of water and saturation, electrical resistivity also
depends on other properties such as porosity, mineralogy, particle shape and orientation and so forth. However, the author believes that the effect of all these properties in all the above tested samples were minimal.

The relationship between friction angle and controlled laboratory samples is shown in Fig. 7. A regression value of $R^2=0.22$ was obtained demonstrating the same trend of increasing friction angle with increasing resistivity. Again, the effect of increase in electrical resistivity and hence friction angle due to decrease in moisture content could be clearly seen from Fig. 8 where the regression value increase to $R^2=0.62$ as a result of variations of moisture content from 25% to 40%. The variation in the compaction blows ranging from 15 to 45 blows resulted in the changes of degree of saturation and could have certainly contributed significantly to the behavior as depicted in Fig. 8.
IV. CONCLUSIONS

In this study, it was observed that the relationship between friction angle and electrical resistivity shows an increasing trend for both the field samples and controlled laboratory samples. Although the points plotted are scattered and the regression values are relatively low, nevertheless, a clear trend of increasing friction angle with increase in electrical resistivity provide an early indication of the possibility to eventually arrive to a better correlation. It is identified here that the behavior displayed is attributed to the increase in moisture content and degree of saturation in all the samples where results of moisture content versus electrical resistivity show strong correlations with significant regression values.

It is suggested here that the framework and focus of the future research should include obtaining enough data covering various types of soils with different moisture content and properties. Correlations should be divided into two main types of soils which are coarse grain and fine grain soils, each having sub divisions for different range of moisture content. Correlations between cohesion and electrical resistivity should also be covered. It is hopeful that with enough data and strong correlations, friction angle ($\phi$) and cohesion ($c$) values could eventually be extracted from the correlated graphs for geotechnical design purpose.

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REFERENCES


Simple and routinely affordable method for therapeutic monitoring of Levetiracetam: A comparison to often applied HPLC method.

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Abstract—Epilepsy is a chronic disease occurring in approximately 1.0% of the world’s population. About 30% of the epileptic patients treated with available antiepileptic drugs (AEDs) continue to have seizures and are considered therapy-resistant or refractory patients. The ultimate goal for the use of AEDs is complete cessation of seizures without side effects. Because of a narrow therapeutic index of AEDs, a complete understanding of its clinical pharmacokinetics is essential for understanding of the pharmacodynamics of these drugs. These drug concentrations in biological fluids serve as surrogate markers and can be used to guide or target drug dosing. Because early studies demonstrated clinical and/or electroencephalographic correlations with serum concentrations of several AEDs, it has been almost 50 years since clinicians started using plasma concentrations of AEDs to optimize pharmacotherapy in patients with epilepsy. Therefore, validated analytical method for concentrations of AEDs in biological fluids is a necessity in order to explore pharmacokinetics, bioequivalence and TDM in various clinical situations. Levetiracetam is a new antiepileptic drug prescribed for the treatment of patients with refractory partial seizures with or without secondary generalization as well as for the treatment of juvenile myoclonic epilepsy. A wide variability in concentration-response relationships is expected in individual patients. Thus, the levetiracetam plasma concentration measurement could be used to help clinicians detect severe intoxication or therapy failure as well as to verify compliance. We applied homogenic enzymatic immunoanalysis for determination of levetiracetam plasma or serum levels in both paediatric and adult neurology patients on maintenance doses of oral drug. We have also compared concentrations measured by high performance liquid chromatography (HPLC) method with ultraviolet (UV) detection as a reference measuring parallelly in an indipendent extramural laboratory with concentrations measured using our method. We found that there is excellent correlation as already published by some authors used other methods including HPLC and UPLC with tandem weight spectrometry. Potentially toxic level is considered as any levels > 37 umol/l, but great concern should be given for existing inter and intra individual variability.

Keywords—AEDs Blood level Determination Methods, EMIT versus HPLC, TDM, Levetiracetam

I. INTRODUCTION

Epilepsy is a common neurological disorder in both paediatric and adult population affecting up to one percent of children for which the mainstay of treatment is anticonvulsant medication. Despite the frequent use of anticonvulsant drugs, remarkably little is known about the safety and efficacy of most of these medications in the paediatric epilepsy population. Many drugs have been used in the past to relieve symptoms by inducing changes in permeability to specific ions thus to stabilize membranes and interfering with release of neurotransmitters or other mechanisms. Although there is continuous emergence of new agents, pharmacoresistant symptom and adverse reactions including drug-drug interactions continue to challenge. Hardly preventable adverse effects (gastrointestinal, mental, or behavioral, neurologic or less commonly cutaneous, haematologic hepatic) by some agents may lead to reduction of quality of life, whereas some agents are known teratogens. Despite all developments in the field, no drug is considered as causal therapy for epilepsy to date. Thus the aim of antiepileptic drug prescription after management of emergency situation like status epilepticus is prophylactic rather than therapeutic. Of 34 anticonvulsants currently approved for use for instance by the US Food and Drug Administration (FDA), only 13 have been approved for use in children being among others since 2012 levetiracetam as an adjunctive treatment for partial onset seizures in infants and children from one month of age as Cornier and Chu concluded that the current data leading to the approval of levetiracetam for use in infants and children with partial onset seizures is encouraging, although more work needs to be done before definitive conclusions can be drawn about the efficacy of levetiracetam across different paediatric age groups[1]. The purpose of this
The company ARK Diagnostics produces the set for analysis of levetiracetam concentration in serum or plasma. The levetiracetam assay is an enzyme multiplied imunoassay technique (EMIT). It’s based on competition between levetiracetam in the specimen and levetiracetam labeled with the enzyme glucose-6-phosphate dehydrogenase (G6PDH) for binding sites of antibodies (rabbit polyclonal antibodies). If the labeled levetiracetam molecules bind to antibodies, enzyme activity is decreased. The activity of G6PDH grows with increasing concentration of the drug from specimen. This enzyme catalyzes conversion of nicotinamide adenine dinucleotide (NAD) to NADH whose concentration is measured spectrophotometrically as a rate of change in absorbance. Samples have been analysed 16 patients by HPLC-UV and EMIT (ARK Diagnostics) and the results were compared.

III. Results

Sixteen samples were measured by EMIT method and within the same time have been also simultaneously determined by HPLC method. Thus, concentrations from all sixteen samples were valid for comparison. There is very good correlation between EMIT and HPLC-UV detection method. Pearson correlation coefficient 0.8265 means strong correlation dependence at 0.05 level of significance. There is also interindividual difference, although the dosing schedules were the same for patients under treatment (Fig. 1).
Table 2. Demonstrates correlation between values measured by compared methods also illustrated Fig. 2. (correlation graph) below.

IV. DISCUSSION

Epilepsy is a chronic disease occurring in approximately 1.0% of the world's population. About 30% of the epileptic patients treated with available antiepileptic drugs (AEDs) continue to have seizures and are considered therapy-resistant or refractory patients. The ultimate goal for the use of AEDs is complete cessation of seizures without side effects. Because of a narrow therapeutic index of AEDs, a complete understanding of its clinical pharmacokinetics is essential for understanding of the pharmacodynamics of these drugs. These drug concentrations in biological fluids serve as surrogate markers and can be used to guide or target drug dosing. Because early studies demonstrated clinical and/or electroencephalographic correlations with serum concentrations of several AEDs, it has been about half a century since clinicians started using plasma concentrations of AEDs to optimize pharmacotherapy in patients with epilepsy. In 1994 it was reported that the cost of uncontrolled epilepsy in the UK was £4167 per patient per year. That of controlled epilepsy was found to be £1630 [2]. If complete seizure control could reduce the costs of managing patients with epilepsy in a significant number of patients, the gains would be appreciable. TDM can help to improve seizure control in numerous ways including:

- identification of therapeutic failure due to under-dosage,
- identification of therapeutic failure in the presence of 'optimal' dosage suggesting that a different AED should be tried,
- detection of non-compliance with prescribed therapy, which may be responsible for unnecessary and avoidable therapeutic failure,
- identification of the uncommon situation in which over-dosage causes increased seizures,
- detection of pharmacokinetic interactions which may compromise the adequacy of the therapy.

Nowadays, therapeutic drug monitoring (TDM) is widely accepted as a method to improve the effectiveness and safety of the first generation of AEDs and to identify an individual's optimum concentration and to individualize drug therapy [3]. Validated analytical methods for determining concentrations of AEDs in biological fluids is a necessity tool both in order to explore pharmacokinetics, bioequivalence studies and therapeutic drug monitoring (TDM). There are abundant published articles on the analysis of specific AEDs by a wide variety of analytical methods in biological samples [4]. A new generation of antiepileptic drugs (AEDs) has reached the market in recent years with ten new compounds: felbamate, gabapentin, lamotrigine, levetiracetam, oxcarbazepine, pregabalin, tiagabine, topiramate, vigabatrin and zonisamide. The newer AEDs in general have more predictable pharmacokinetics than older AEDs such as phenytoin, carbamazepine and valproic acid (valproate sodium), which have a pronounced inter-individual variability in their pharmacokinetics and a narrow therapeutic range. For these older drugs it has been common practice to adjust the dosage to achieve a serum drug concentration within a predefined ‘therapeutic range’, representing an interval where most patients are expected to show an optimal response. However, such ranges must be interpreted with caution, since many patients are optimally treated when they have serum concentrations below or above the suggested range. It is often said that there is less need for therapeutic drug monitoring (TDM) with the newer AEDs, although this is partially based on the lack of documented correlation between serum concentration and drug effects. Nevertheless, TDM may be useful despite the shortcomings of existing therapeutic ranges, by utilisation of the concept of ‘individual reference concentrations’ based on intra-individual comparisons of drug serum concentrations. With this concept, TDM may be indicated regardless of the existence or lack of a well-defined therapeutic range. From the ten newer AEDs, all have different pharmacological properties, and therefore, the usefulness of TDM for these drugs has to be assessed individually. For vigabatrin, a clear relationship between drug concentration and clinical effect cannot be expected because of its unique mode of action. Therefore, TDM of vigabatrin is mainly to check compliance. The mode of action of the other new AEDs would not preclude the applicability of TDM. For the prodrug oxcarbazepine, TDM is also useful, since the active metabolite carbazepine is measured. For drugs that are eliminated renally completely unchanged (gabapentin, pregabalin and vigabatrin) or...
mainly unchanged (levetiracetam and topiramate), the pharmacokinetic variability is less pronounced and more predictable. However, the dose-dependent absorption of gabapentin increases its pharmacokinetic variability. Drug interactions can affect topiramate concentrations markedly, and individual factors such as age, pregnancy and renal function will contribute to the pharmacokinetic variability of all renally eliminated AEDs. For those of the newer AEDs that are metabolised (felbamate, lamotrigine, oxcarbazepine, tiagabine and zonisamide), pharmacokinetic variability is just as relevant as for many of the older AEDs. Therefore, TDM is likely to be useful in many clinical settings for the newer AEDs. The purpose of the present review is to discuss individually the potential value of TDM of these newer AEDs, with emphasis on pharmacokinetic variability. [5].

The rationale for the determination of AEDs and their metabolites in body fluids and tissues arises from different fields of investigations and clinical situations. Either drug or metabolites levels are required for regular monitoring of therapeutic drug levels, for adverse drug reactions, for drug-drug interaction studies, for issues of toxicity concern, for pharmacokinetic, pharmacokinetic/pharmacodynamic and bioequivalence studies. AEDs are often used in polypharmacy including up to three different AEDs, each of them having several own metabolites [6]. TDM is more important for drugs with a narrow therapeutic range, where a correlation has been established between drug concentration and its therapeutic and toxic effects. Although reasonably well-defined target ranges in serum concentrations have been determined for most of the established AEDs [7, 8], it should be remembered that these ranges only became established after the development and general availability of sensitive and reliable analytical methods. Thus, even though phenobarbital and phenytoin became use for clinical application in the early 1900s, after the development of analytical methods in the 1960s, it was only since the early 1970s that target ranges were identified [9, 10]. Since then, monitoring AEDs such as carbamazepine, vaproate and ethosuximide has also become widely accepted in clinical practice [11]. Although the target ranges have been defined for some of the AEDs, the true therapeutic range is defined for a given patient as the concentration that prevents occurrence of epileptic episodes without causing side effects. Since 1989, several AEDs have been have been approved for clinical use, and because their regulatory trials were not serum concentration controlled or designed to investigate the relationship drug concentration and effect, the value of monitoring these drugs is presently controversial. However, some of the newer AEDs have pharmacological properties suggesting that their optimal use may be facilitated by use of TDM, and this has been the subject of recent valuable debate [12, 13]. The AEDs have been measured by a wide variety of analytical methods in serum, blood, saliva, urine and tissue. For the classic AEDs (carbamazepine, ethosuximide, phenobarbital, vaproate) and some of the new AEDs (felbamate, topiramate, zonisamide etc), automated enzyme-multiplied immunoassay technique (EMIT) and fluorescence polarization immunoassays (FPIA) are available and allow rapid and accurate determination of concentrations in biological fluids, usually serum or plasma. For other AEDs laboratories rely on chromatographic methods; gas-chromatography (GC) and high-performance liquid chromatography (HPLC) with a various detection methods, which are more labor-intensive and relatively more expensive. A number of simultaneous chromatographic assays for AEDs have been developed in the past. The early initial simultaneous AED assays, from 1970s and 1980s, concentrated on separating the older AEDs such as ethosuximide, primidone, carbamazepine, carbamazepine-10,11-epoxide, phenytoin, and phenobarbital ([14-18] Many subsequent assays separated the same compounds with the inclusion or removal of one or more additional drugs or metabolites such as ethylpenacemide [16], 5-parahydroxyphenyl-5-phenylhydatoin, N-desmethylnmethsuximide [17], and lamotrigine and phenyl-2-thyl-malonomamide [19]. Another inclusive assay separated ethosuximide, primidone, carbamazepine, phenytoin, Phenobarbital, carbamazepine metabolites, phenobarbital metabolite and felbamate [20]. All of these assays employ ultraviolet (UV) detection, thereby increasing the risk of metabolite or metrix interferences. Assay developed over past 15 years have focused more separating newer AEDs. There are also new technological advance in the use of capillary electrophoresis (CE) for TDM. Like other chromatographic methods, CE allows simultaneous measurement of several AEDs and can provide automation of procedures, low cost, and rapid speed with high specificity [21, 22]. The number of articles with the analytical assay of AEDs in various biological matrices is increasing in accordance with growing interest in the situations with pharmacokinetic, TDM and bioequivalence studies of clinical and research fields. The present review was to focus to current technologies applied to the analysis of AEDs in biological media for monitoring individual AEDs or simultaneous monitoring of AEDs in recent years. Levetiracetam is among the latest AEDs to be licensed for clinical use indicated for adjunctive treatment of intractable partial seizures that do not respond to other AEDs. Its mechanism of action is not clearly defined, [23], whereas its pharmacokinetics such as rapidly absorption from the gastrointestinal tract, non-significant protein binding and elimination of linear pharmacokinetics is well kown It is minimally metabolised and excreted essentially unchanged via the kidneys. Levetiracetam has no propensity to interact with other AEDs [24].There is little published information on levetiracetam blood levels in patients with epilepsy. Nevertheless, audit of the clinical trials data for levetiracetam suggests a target range of 35 - 110 µmol/L, while blood levels can be measured by HPLC with ultraviolet detection [25]. Levetiracetam is a new antiepileptic drug prescribed for the treatment of patients with refractory partial seizures with or without secondary generalization as well as for the treatment of juvenile myoclonic epilepsy. A rapid and specific method by high-performance liquid chromatography diode array detection was developed to measure the concentration of levetiracetam in human plasma. A wide variability in concentration-response relationships was observed in patients.
Nevertheless, the levetiracetam plasma concentration could be used to help clinicians detect severe intoxication or to verify compliance by repeating the measurement in patients [26]. An isocratic high performance liquid chromatographic micromethod, Ratnaraj et al. reported that no interference from commonly prescribed antiepileptic drugs (carbamazepine and its metabolite carbamazepine epoxide, ethosuximide, gabapentin, lamotrigine, phenobarbital, phenytoin, primidone, valproic acid, and vigabatrin) was observed, and thus the method can be used to monitor levetiracetam in patients on polytherapy antiepileptic drug regimens. [27]. In a study to evaluate the efficacy and tolerability of levetiracetam (LEV) as add-on therapy in children with refractory epilepsies and to determine the value of LEV blood level monitoring in this population, Giroux et al published that the most frequently observed adverse effects were drowsiness, behavioral difficulties, increase in seizure frequency and headaches. The majority (60.5%) of the responders received doses between 10 and 50mg/kg/day and had a plasma concentration between 5 and 40microg/ml, however, no clear correlation between drug plasma concentration and efficacy has been found. [28]. According to the study to evaluate the efficacy and tolerability of Levetiracetam in a large pediatric cohort with drug-resistant epilepsy from a prospective multicenter observational study it has been reported that Levetiracetam is a well-tolerated new antiepileptic drug that may effectively improve seizure control as an add-on drug in resistant epilepsy in childhood with good tolerability, however, neurologically handicapped children appear at increased risk for reversible neurocognitive side effects and have a poorer treatment response. [29] Levetiracetam is among the most recently licensed antiepileptic drug (AED) for adjunctive therapy of partial seizures. Its mechanism of action is uncertain but it exhibits a unique profile of anticonvulsant activity in models of chronic epilepsy. [30]. Most studies suggest that levetiracetam is effective against partial and generalized epilepsy. In resistant partial epilepsy, the percentage of responders reaches 64%, with 8 to 23% seizure free. Levetiracetam is used to treat symptomatic and idiopathic epilepsies. The drug has also proven effective against photosensitivity and epileptic and nonepileptic myoclonus. The most frequent side effects involve the behavioral sphere and manifest mostly in patients with a history of behavioral problems. In some patients, levetiracetam increases the number of seizures, but this adverse reaction can be partially avoided with slow titration. Doses for children should be 130 to 140% of those advised for adults. [31]. Treatment of seizures in pediatric patients is complicated by the fact that the etiology of the disorder and the pharmacokinetics, efficacy, and safety of antiepileptic drugs may differ from that in adults. With few controlled clinical trials of AEDs in children, the selection of agents to treat paediatric patients must be made on the basis of information from small uncontrolled studies or the extrapolation of clinical trial results in adults. Data from a large number of children with a wide range of seizure disorders who were treated in small-scale prospective studies, or whose records were retrospectively evaluated, indicate that levetiracetam reduces the frequency of seizures in pediatric patients. Available data also indicate that levetiracetam is well tolerated in pediatric patients, with a safety profile similar to that in adults, a low potential for behavioral disturbances, and no reported idiosyncratic adverse reactions. As with other AEDs, children metabolize and clear levetiracetam more rapidly than adults, and somewhat higher doses (based on body weight) are needed to achieve desired plasma concentrations. Several ongoing studies will provide further information on the pharmacokinetics, efficacy, and safety of levetiracetam in this patient population.[32]. Some study results support the use of a weight-based LEV dosing regimen and provide a basis for a recommended pediatric dosage regimen, but the relationship between LEV plasma concentrations and clinical effect has not been evaluated fully and could differ between adults and children and Clinical studies should be able to validate these dosing recommendations. [33]. Some significant covariates for pharmacokinetics of levetiracetam were identified: (a) age on the absorption rate constant ; (b) bodyweight, dose, clearance, and concomitant enzyme-inducing AED on plasma oral clearance (CL/F); and (c) bodyweight on the apparent volume of distribution after oral administration (V(d)/F). The main explanatory covariates were age on absorption rate constant, bodyweight on CL/F and V(d)/F, and enzyme-inducing AED on CL/F, of which bodyweight was the most influential covariate. However, the most influential covariate of levetiracetam pharmacokinetics in children is bodyweight. [34]. Both serum and CSF levetiracetam concentrations rose essentially linearly and dose-dependently, suggesting that transport across the blood-brain barrier is not rate limiting over the levetiracetam concentration range observed. However, while apparent elimination half-life (t1/2) values for both serum and CSF were dose-independent (mean value range 1.8-2.8 and 4.4-4.9 h, respectively), t1/2 values for CSF were significantly larger. As the serum free/total serum levetiracetam concentration ratio (free fraction) was 1.01 +/-0.02 (mean +/- S.E.M.), it can be concluded that levetiracetam is not protein bound. [35]. Levetiracetam is completely and rapidly (Tmax, 1 h) absorbed after oral ingestion with bioavailability of 100% [36, 37]. Although drug-food interactions do not show on the extent of absorption, rate of absorption is slowed in the presence of food. Administration of a crushed LVT tablet together with 120 ml of an enteral nutrition formula has been associated with a mean 27% decrease in peak LVT concentration, but the effect was not statistically significant. [38]. Levetiracetam shows linear pharmacokinetic and renal elimination with approximately 66% of a dose eliminated unchanged and 27% as inactive metabolite [39]. The relationship between LVT serum concentrations and clinical effect has not been ascertained, and consequently the value of serum concentrations measurements is not established. Because of its favorable therapeutic index, low plasma protein binding and minimal side-effect profile, routine monitoring of LVT serum concentrations appears to be unnecessary for safe use of the drug, and dosing can be readily guided by the therapeutic response [40, 41]. Nevertheless, its use in ascertaining
compliance and managing patients that are overdosed would be helpful. Because LVT has a relatively short half-life, sampling time in relation to dose ingestion is important for the interpretation of the drug concentration. Ideally samples for LVT measurements should be drawn before the morning dose. Because LVT can undergo in vitro hydrolysis, it is important to separate whole blood from serum as soon as possible so as to avoid LVT hydrolysis that would result in spuriously lower concentrations being measured [42]. Although AED monitoring in saliva may some clinical applicability, it has not yet come into routine use, but, a significant positive correlation exists between LVT saliva and serum concentrations, LVT like other AEDs, can be measures in saliva as an alternative to blood-based assays for minotoring the LVT therapy [43, 44]. Numerous chromatographic methods have been reported for the quantification of LVT in biological fluids. Different HPLC methods for the determination of LVT in human plasma have been reported coupled with UV [45, 46], or diode array detection [47, 48]. mostly after sample pretreatment by expensive solid-phase extraction or time-consuming liquid-liquid extraction procedures [46, 48]. The availability of simple, accurate and inexpensive analytical assays is crucial for the successful use of TDM in clinical practice. LVT spiked plasma sample preparation by different kinds of deproteinization before HPLC-diode array detection was first explored by Pucci et al[47], and subsequently applied to patient samples analysis by HPLC-UV [49, 45, 50]. Otherwise, these involve GC-NPD [51] and GC-MS [52]. Most of the reported methods lack selectivity, sensitivity, and reliability. They encounter also problems particularly tedious and time-consuming sample preparation as well as high sample volume. Recently, LC-tandem MS is considered a gold standard to utilize in analysis of drugs in biological fluids. The high sample throughput, selectivity and sensitivity for analytes of interest may increase the applicability of tandem MS in clinical chemistry as well as clinical studies [53-55].

V. CONCLUSIONS

We found that there is excellent correlation as already published by some authors used other methods including HPLC and UPLC with tandem weight spectrometry. Potentially toxic level is considered as any levels > 37 umol/l, but great concern should be given for existing inter and intra individual variability. Our method is suitable for routins levetriracetam plasma level determination in the process of therapeutic drug monitoring.

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New tree species for agroforestry and energy purposes

Andrea Vityi, Béla Marosvölgyi

Abstract—Bibliographic data and the results of present domestic experiments show that Paulownia species can be grown easily, have high biomass production, favourable energetic parameters and modest requirements as to site quality. Establishment of Paulownia plantations may support the aim to meet the growing needs for site-remediation and biomass for energy purposes. This paper gives an outlook on the relevant research activity and results of the University of West Hungary. The use of the selected species can reasonably be recommended for these purposes, since the plantations - grown from self-developed propagating material, managed and tested according to the formulated methods - have a high production of biomass and thus confirmed as suitable for energy purposes. The use of Paulownia for agroforestry purposes may also be a prospective way of multifunctional landuse while providing beneficial ecosystem services.

Keywords—energy, biomass, Paulownia, agroforestry.

I. INTRODUCTION

The empress tree (Paulownia tomentosa) is one of the world’s most multifaceted tree species. It originates from China and are also grown throughout Asia, USA, Australia and Europe. (AFBI, 2008) Its wide spectrum of utilization ranges from industrial use (furniture, timber, pulp and paper, energy), to medical and hive products, or decoration (eg. ornamental trees and wood carvings).

In recent years the awareness and business interest for Paulownia is growing fast in Central-Eastern Europe, while there is a lack of research activities in the subject of cultivation, utilization, and adoption of the best practices in these countries. However it is of utmost importance to take the local conditions into consideration when adopting new species and technologies.

As no preliminary research activity can be recognized on this subject in Hungary, NyME KKK together with external partners (cooperative and farmers) started a research program on energy and agroforestry use of Paulownia.

II. MATERIAL AND METHOD

According to the literature, Paulownia can be multiplied from seeds, cuttings or by micropropagation. (Al-Tinawi, I. A. et al., 2010, Lobona, S., 2008) (Gyuleva, V., 2008) In favourable conditions a 10 year-old tree may reach 30-40 cm in diameter, 10-12 m in height and provide 0,2-0,6 m³ volume production. (Yang, 2004)

In the first stage of our experiments we used selected mother plants and developed a special method of propagation form seeds. The method is based on a special substrate composition and strict planting protocol. In spring seedlings with 4-6 leaves were relocated into planting containers where they grew until planting out in autumn.

The first experimental bioenergy plantations were established from this selected material in autumn 2006. Planting and research activities were coordinated by the Eco-energetic Research Division of NyME KKK.

Plantations were located in various parts of the country with different climate conditions. Also different planting structures and cutting rotation systems were applied in order to examine the growing and other relevant parameters of the plants.

In order to scan the usability of the trees in energetic processes we studied the international literature available on Paulownia, then made tests with the samples originated from our own plantations on bulk density, moisture content, ash content, and heating value, which are basic parameters concerning energetic utilization.
III. RESULTS AND CONCLUSIONS

A. Experiences on propagation, planting and plant management

By the use of specifically selected Paulownia plants and self-developed propagation method we managed to realize 80% plant survival. Further experimental plantations based on the selected material are planned to be established for extended research purposes in the next years.

B. Experimental results of biomass production

According to the results of the crop yield survey the biomass production of the experimental plantation was definitely high (55 t/ha).

It has to be underlined that the given high yield volume is only valid under the specific parameters of the experimental system (applied planting structure, site conditions, selected material, cutting rotation, site management, etc.). Furthermore, given the dimensions of the test parcel, we had to calculate with border-effect which surely had significant benefits for the biomass production volume.

C. Results of the energetic analyses of the biomass form the Paulownia experimental site

Having studied the literature we found that the available data on certain energetic parameters of Paulownia-biomass varied, but basically were comparable with the results of our measurements. (Table 1) (Source: Vityi-Marosvölgyi, 2012)

The good test results of the woody biomass from the experimental sites - whether in annual cutting rotation or in multi-year rotation system - show that Paulownia may definitely be suitable for energy purposes.
### Table 1 Results of the University of West Hungary Cooperatorial Research Centre’s tests compared with the data available in the literature on the energetic parameters of Paulownia

<table>
<thead>
<tr>
<th></th>
<th>Mechanical and combustion parameters of the one-year-old shoot of selected Paulownia tomentosa variety</th>
<th>Bibliographical data on the mechanical and combustion parameters of Paulownia species</th>
</tr>
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<tbody>
<tr>
<td></td>
<td>wood</td>
<td>bark</td>
</tr>
<tr>
<td>Bulk density (g/cm³, air-dry)</td>
<td>0,35</td>
<td></td>
</tr>
<tr>
<td>Moisture content (m/m%, air-dry)</td>
<td>11,80</td>
<td>10,90</td>
</tr>
<tr>
<td>Ash content (m/m%)</td>
<td>0,92</td>
<td>3,05</td>
</tr>
<tr>
<td>Heating value (MJ/kg, air-dry)</td>
<td>16,66</td>
<td>17,40</td>
</tr>
<tr>
<td>Heating value (MJ/kg, absolute dry)</td>
<td>18,92</td>
<td>19,67</td>
</tr>
</tbody>
</table>

*based on bark-rate of 21% measured from one-year-old shoot samples

Based on the positive bibliographic data on the use of Paulownia in alley-cropping and the favorable results of NyME KKK’s field tests we decided to extend the examinations to the use of Paulownia for agroforestry purposes.

Agroforestry experiments with Paulownia started in 2012. The initial step was the plantation of the first and so far only experimental Paulownia intercropping system in Hungary. In the next years NyME KKK plan to make investigations on crucial parameters of sustainable management of the Paulownia-intercrop systems. This activity will also be a contribution to EU FP7 project called „AGFORWARD“. This is a four-year pan-european project, with the goal of promoting appropriate agroforestry practices that advance sustainable rural development.

### REFERENCES


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**D. The Use of Paulownia in agroforestry systems**

Agroforestry is the practice of deliberately combining woody vegetation (trees and/or shrubs) with crop and/or livestock systems.

Agroforestry practices help farmers to diversify farm income, while benefiting from the resulting ecological and economical interactions eg. improve soil and water quality, reduce erosion, pollution, or damage due to extreme weather, enhance resource efficiency, biodiversity, and the resiliency of the production system. Agroforestry systems manifest in several practices: forest buffers, windbreaks, silvopasture, alley cropping, forest farming, etc.
Extremely delayed elimination of methotrexate in a young man with osteosarcoma: A case study demonstrating an association with impaired renal function.

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Abstract— Background: Methotrexate (MTX) is one of the most widely used anti-cancer agents, and administration of high-dose methotrexate followed by leucovorin (LV) rescue therapy is an important component in the treatment of a variety of cancers. High-dose MTX is thought to be safe for administration in patients with normal renal function with concomitant alkalinization, hydration, and pharmacokinetically-guided leucovorin rescue. However, acute renal failure and other adverse outcomes are unavoidable under certain circumstances. High-dose methotrexate (HDMTX)-induced renal dysfunction can be a life threatening event, because it delays methotrexate excretion, thereby exacerbating the other toxicities of MTX. High-dose methotrexate -induced renal dysfunction continues to occur in patients with osteosarcoma who are treated on clinical protocols despite optimal supportive care. Approximately, 1.8% of patients with osteosarcoma, who received HDMTX may develop nephrotoxicity of various grade carrying considerable mortality rate. High-dose methotrexate-induced renal dysfunction can be life threatening, because it delays methotrexate excretion, thereby exacerbating the other toxicities of methotrexate. HDMTX-induced nephrotoxicity has been usually managed with high-dose leucovorin, dialysis-based methods of MTX removal, thymidine, and recently by administration of the recombinant enzyme, carboxypeptidase-G2 (CPDG2), which cleaves MTX to inactive metabolites. Objective: The aim of this paper is to describe the case of an adult Caucasian male patient with osteosarcoma who presented with extremely delayed MTX clearance after high-dose administration conducted according to the EURAMOS protocol. In the present case, we discuss also the fate of the patient where delayed MTX excretion was a big challenge and finally managed using supportive measures including high doses of leucovorin and very effectively CPDG2 without having to adopt other invasive procedures like dialysis.

Keywords—Osteosarcoma, High-doseMXT, Nephro-Toxicity

I. INTRODUCTION

Osteosarcoma (osteogenic sarcoma) is the most common type of primary bone cancer most occurring in children and young adults with peak incidences in adolescence and at age >60 years, but can occur at any age.[1] Methotrexate (MTX), a
classic antifolate is one of the most widely used and well studied anticancer agents, where the administration of MTX doses $\geq 1000$ mg/m$^2$ combined with leucovorin (LV) rescue is defined as high-dose methotrexate (HDMTX) and is an important component of treatment for a variety of malignancies, including osteosarcoma. Todate, a very high-dose methotrexate usually $> 1$ g/m$^2$ administered as an intravenous infusion remains an important component in the treatment of variety of cancers in particular for osteosarcoma, but this HDMTX treatment schedule carries a risk of nephrotoxicity among others. Data from a number of studies performed in the 1970s showed that a sustained elevation of serum MTX concentrations at 24 h ($\geq 5$ μmol/L), 48 h ($\geq 1$ μmol/L) and 72 h ($\geq 0.1$ μmol/L) after the start of the MTX infusion is considered to be toxic referring to the usual serum MTX level $< 0.1$ μmol/L 48 h after HDMTX administration. In the era of optimal supportive care the incidence of grade 3–4 ARF after HDMTX has reportedly decreased markedly in solid-cancer patients. Although HDMTX-associated severe acute renal failure (ARF) is an infrequent, the very high doses of MTX generally used in solid-cancer patients such as with osteosarcoma, who receive a MTX dose $> 8$ g/m$^2$, the risk of severe ARF may be certainly higher. According to large case series, the reason why an individual patient comes prone to develop ARF after HDMTX despite modern supported care remains unexplained in the majority of cases (1-3). However, drug drug interactions are mostly unrecognized or overlooked. In case of our patient the role of pre treatment with Cis-platin in the near past before the initial therapy with MXT may be the partial explanation for the precipitation of acute renal dysfunction grade 2-3 (Renal toxicity was graded using World Health Organization criteria (Grade 1, serum creatinine levels $< 1.5 \times$ ULN; Grade 2, 1.5–3.0 × ULN; Grade 3, 3.1–6.0 × ULN; and Grade 4, $> 6.0 \times$ ULN) Although the incidence and mortality of HDMTX-induced renal dysfunction appear to have decreased significantly since the 1970s, nephrotoxicity continues to occur and may be fatal. Therefore, in situations, where usual care fails in patients with delayed MTX excretion and high plasma MTX concentrations, other measures like treatment by CPDG$_2$ should be considered to lower plasma MTX concentrations rapidly and efficiently as previously recommended. The median time to recovery of renal function for a number of patients, as defined by the individual studies, was 16 days (range, 4–48 days), where treatment of patients did not include carboxypeptidase-G$_2$. According to published literature reviewed from 1977 to 2002 on recovery of renal function in patients with methotrexate-induced renal dysfunction, the median time to recovery of renal function for a number of patients, as defined by the individual studies, was 16 days (range, 4–48 days), where treatment of patients did not include carboxypeptidase-G$_2$. [3], [8]

II. CASE DESCRIPTION

A 37-year-old Caucasian male had been treated initially with a combination of doxorubicin and cisplatin for proven diagnosis of osteosarcoma. Just after a month later, the patient was scheduled for methotrexate treatment according to the EURAMOS (European and American Osteosarcoma Study Group) protocol, a joint protocol of four of the world’s leading multi-institutional osteosarcoma groups; which uses a dose of 12 g/m$^2$ over a 4-hour infusion and repeated with 11.34 g/m$^2$. The concentration of MXT determined by fluorescent polarization immunoassay (FPIA) method 6 hours post-infusion showed levels within expected range. However,
measurements taken 24 hours post-dose and later were extremely high. This indicated poor elimination and was also confirmed by significantly elevated serum creatinine (Fig. 1) as well as blood urea nitrogen. Drug plasma level monitoring was continued on a daily basis as per protocol guidelines until the level reached less than 0.1µM. This took one month +8 days from initial MXT administration.

Fig. 1. Extremely delayed methotrexate (MXT) elimination corresponding with serum creatinine levels (S-Cr) and persistence of high blood urea nitrogen (BUN) demonstrating significant renal function impairment.

Taking into consideration that with high-dose methotrexate, toxic concentrations are generally considered to be: ≥ 5 µmol/L at 24 hours after the dose, ≥ 0.5 µmol/L at 48 hours, and ≥ 0.055 µmol/L at 72 hours; we declare our findings as potentially extremely toxic levels. The test results are used to guide the amount and timing of leucovorin (folinic acid) given as a "rescue" treatment, but the effect of the rescue therapy was not satisfactory in this case. Finally, carboxypeptidase-G2 (CPDG2) has been used with significant effect in reducing the drug level by 80 % of the previously recorded value. However, renal function and further drug elimination were lagging for several days. BUN and serum creatinine were not restored to normal until after a month. It took more than one month to get the drug plasma level of 0.11 umol/l as illustrated above (Fig.1). At the time of very high drug levels, the patient also manifested with significantly high activity of liver aminotransferases, namely ALT (up to 30 U/L). AST was moderately increased (4 U/L) and both were shortly thereafter restored. This was in contrast to BUN and serum creatinine levels which remained abnormal over the course of a month. Thrombolytic and leukocyte profiles were also demonstrably unstable throughout follow-up until the elimination of the drug (Fig. 2). Prominent leukopenia was observed in the week after drug exposure; whereas thrombocytopenia was a few days earlier (Fig. 2). Both events of leukopenia and thrombocytopenia had several phases demonstrating instability of the blood count in association with a prolonged elevated level of methotrexate exposure.

Fig. 2. Thrombolytic and leukocyte profiles demonstrating instability of the blood count associated with a prolonged elevated level of methotrexate exposure.

### III. DISCUSSION

Some acute toxicities like liver toxicity manifesting with ALT and AST transient elevation were reported in the past as reversible without further concern as self limited [9] as also has been observed in our patient case. Despite advanced management and care measures, high-dose MTX-induced renal dysfunction continues to occur in approximately 2% of patients with osteosarcoma treated in clinical trials. [3] Early recognition and treatment of MTX-induced renal dysfunction are essential in preventing potentially life-threatening toxicities; especially myelosuppression, renal failure, mucositis and dermatitis. In addition to conventional treatment approaches, dialysis-based methods have been used to remove MTX with limited effectiveness. More recently, CPDG2, a recombinant bacterial enzyme that rapidly hydrolyzes MTX to inactive metabolites, has become available for the treatment of MTX-induced renal dysfunction. [4] Certain circumstances like ascites and packed red blood cell infusion may
function as a reservoir and enhance prolonged high level exposure to methotrexate during a high-dose regimen, but our patient had only suffered mild pleural effusion, not ascites, to serve as a possible reservoir. Some previously published studies also identified several clinical variables that influence MTX disposition that, when modified, can reduce the frequency of high-risk MTX concentrations and toxicity [10]. Clearance is exceptionally variable in individuals and association with age and gender has been also documented [11]. However, none of these variables explain the extremely delayed elimination of the drug in our patient. Similar to other antimetabolites, critical determinants of MTX cytotoxicity is not only drug concentration, but also the duration of exposure. High concentrations of MTX may be well-tolerated for brief periods of time; whereas prolonged exposure to low concentrations can result in life-threatening toxicity. The type of toxicity observed with MTX is also a function of this concentration–time dependence. Exposure to millimolar concentrations of MTX for minutes to hours may lead to acute renal, central nervous system, and liver toxicity. Exposure to MTX concentrations as low as 0.01 and 0.005 μM for > 24 hours may result in bone marrow and gastrointestinal epithelial toxicity, respectively [12]. The MTX-induced renal dysfunction is believed to be mediated by the precipitation of MTX and its metabolites in the renal tubules [13-15] or via a direct toxic effect of MTX on the renal tubules [16]. Urinary NAG:creatinine ratio in our patient after 3 weeks continued to demonstrate abnormality correlating to delayed function reversibility since more than 90% of MTX is cleared by the kidneys [17]. MTX is poorly soluble at acidic pH and its metabolites, 7-OH-MTX and DAMPA, are six- to tenfold less soluble than MTX [13, 18]. An increase in urine pH from 6.0 to 7.0 results in a five- to eightfold greater solubility of MTX and its metabolites; a finding that underlies the recommendation of i.v. hydration (2.5–3.5 litres of fluid per m² per 24 hours, beginning 12 hours before MTX infusion and continuing for 24–48 hours) and urine alkalinization (40–50 mEq sodium bicarbonate per liter of i.v. fluid prior to, during, and after the administration of high-dose MTX as performed in the present case. Several drugs have also been associated with increased toxicity when co-administered with MTX.[6] The most significant interactions involve agents that interfere with MTX excretion, primarily by competing for renal tubular secretion, such as: probenecid, salicylates, sulfisoxazole, penicillins, and nonsteroidal anti-inflammatory agents [19], but all were excluded in the present case. MTX-induced renal dysfunction results in sustained, elevated plasma MTX concentrations; which in turn may lead to ineffective rescue by leucovorin and a marked enhancement of MTX’s other toxicities; especially myelosuppression, mucositis, hepatitis, and dermatitis [20, 21]. Previous studies demonstrated that: (a) sustained elevation of plasma MTX concentrations at 24 hours (> 5–10 μM), 48 hours (> 1.0 μM), and 72 hours (> 0.1 μM) after administration of MTX are predictive for the development of toxicity; (b) in the absence of elevated plasma MTX concentrations, the risk for the development of MTX-associated toxicities is minimal; (c) in most circumstances, the development of MTX-associated toxicities can be ameliorated or prevented when patients with elevated plasma MTX concentrations receive pharmacokinetically-guided doses of LV rescue. These studies resulted in uniform institution of aggressive hydration, alkalinization, and pharmacokinetically-guided LV. Nomograms guiding the duration and degree of rescue with LV based upon plasma MTX concentrations as a function of time of drug administration were developed and are being used in ongoing clinical trials that administer high-dose methotrexate [22]. Goren et al., [23] found that analysis of the changes in these sensitive markers of renal tubular damage permitted detection of subclinical methotrexate-induced nephrotoxicity. According to the authors, persistent rises in NAG as well as increased serum creatinine levels in patients with osteogenic sarcoma who were receiving combination chemotherapy that included 12 doses of methotrexate (12 g/m²) was associated with doses of methotrexate that followed the administration of cisplatin (400 mg/m²), while the biphasic pattern of NAG excretion observed in patients suggests more than one mechanism of methotrexate-induced nephrotoxicity. Thus, monitoring renal tubular damage in patients who are receiving methotrexate in combined drug regimens would provide useful information. In the study to determine the risk of impaired excretion of methotrexate in patients with osteosarcoma, who were also receiving cisplatin, it has been found that MTX clearance was impaired in patients with urinary NAG concentrations greater than 1.5 U/mmol creatinine or greater than 50% increase in serum creatinine relative to the pre-therapy level, were approximately 30 times more likely to have MTX half-lives greater than 3.5 hours than were patients with lower values for these markers (i.e., MTX clearance was always impaired if both markers were elevated) [24]. These findings demonstrate that urinary NAG and serum creatinine levels, measured before MTX administration, can be used to identify patients who will have difficulty clearing the drug and thus can be used to guide
rescue measures including: highly effective carboxypetidase-G2, thymidine, and leucovorin in patients at high risk for developing life-threatening methotrexate toxicity after the onset of methotrexate-induced nephrotoxicity and delayed methotrexate excretion [25]. In the case described here, the NAG:creatinine ratio was abnormal several days post-MXT exposure. Moderate hypokalemia and hyponatremia observed later on may also be explained by poor tubular function. Over all Early recognition of the failure to eliminate excess MTX may help start of more effective rescue treatment such as CPDG2 application recently published as better alternative.[5] Based on our literature review and information available from other databases [3], approximately 1.8% of patients with osteosarcoma who are treated with HDMTX develop significant nephrotoxicity at some time during treatment, and the mortality among these patients is estimated at 4.4%. Limitations of this estimate include the following: nephrotoxicity was not defined uniformly across studies, trials were not designed to capture data concerning nephrotoxicity comprehensively, and few published studies included a comprehensive description of MTX-related toxicities. In addition, only patients entered on clinical trials, who may be more likely to receive optimal supportive care, were included in the estimate. Therefore, the incidence of MTX-induced renal dysfunction among all patients receiving HDMTX may be higher.

Myelosuppression occurs in 28 per cent of the patients and in 8 per cent of the courses and usually results from delayed MTX excretion secondary to mild reversible nephrotoxicity.[26] However, the manifestation in a present case osteosarcoma patient was not relevant for significant haematotoxicity. Aggressive hydravion, urine alkalinization aiming to keep urine pH 7 to 8.5, Leucovorin administration, close monitoring of urine output, fluid balance avoiding negative balance, serial MTX levels and Serum Cr. monitoring are among preventive measures against HDMTX potential nephrotoxicity. [27] In some individuals preventive measures sometimes fail to protect patients from MTX-induced nephrotoxicity. In such cases, standard approaches namely, hydration, urine alkalinization, and leucovorin rescue are also key to managing patients who develop HDMTX-induced renal dysfunction. Additional measures include extracorporeal interventions such as hemofiltration and dialysis, as well as pharmacotherapy with glucarpidase. Glucarpidase (carboxyptidase G2), which was approved by the US Food and Drug Administration (FDA) in January 2012 for the treatment of cases with plasma MTX concentrations (> 1 μM) in patients with delayed MTX clearance due to impaired kidney function despite standard leucovorin application. [28]

IV. Conclusion

Based on the case demonstrated, We recommend prompt recognition of patients with poor elimination of after administration of high-dose methotrexate and to start effective rescue therapy including glucarpidase (carboxyptidase G2) for its effective elimination of MTX to avoid further deterioration of health and to improve overall outcomes.

V. Acknowledgements

Conflict of interest disclosure: The corresponding author declares no conflict of interest pertaining to this case report.

References

Transmission Corridor between Romania-Moldova-Ukraine

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Abstract – Constructed in 1986, the 750 kV line connecting the Ukrainian and Romanian transmission networks went out of service in the mid-1990s due to damage to the lines. Although the Romanian TSO (Transelectrica) and the Ukrainian TSO (Ukrenergo) carry plans to restore the line, each has experienced significant development of their transmission networks since the line went out of service. This article identifies the optimal configuration of the corridor to serve the transmission requirements of the system operators in Romania, Ukraine and Moldova. Currently the transmission corridor, which had consisted of a 750kV AC Over Head Line (OHL), is not in operation and is in a state that cannot be easily repaired. The OHL has been damaged so that it could be considered as “non-existent” for each party. The investment scenarios themselves are comprised two voltage levels considered for the corridor: 400 kV and 750 kV. In turn, these voltages can be analyzed in terms of synchronous AC or asynchronous DC connection via a back-to-back station that may be located in either Moldova or Romania.

Keywords: asynchronous, back-to-back, IPS/UPS, RUM

I. INTRODUCTION

Currently the Romanian – Moldovan – Ukrainian (RUM) transmission corridor, which had consisted of a 750kV AC Over Head Line (OHL), is not in operation and is in a state that cannot be easily repaired. The OHL has been damaged so that it could be considered as “non-existent” for each RUM party. The existing route of the old 750 kV transmission line is depicted in Fig.1.

Fig. 1. The route of the old 750 kV transmission line

Although the original transmission corridor is directly between Pivdennoukrainska NPP (Ukraine) and Isaccea (ROM) Substations, Ukraine plans to construct a new 750kV OHL between Pivdennoukrainska NPP and Primorska Substation (see red dashed line in Fig.2). Hence, the corridor under discussion in this article will cover the existing Right of Way between Primorska (Ukraine)-Isaccea (ROM) substations. This corridor is depicted in Fig.2 (blue dashed line).

The transmission line distances between the substations in the corridor are provided in Table 1. The transmission lengths do not represent fly-over distances but rather the total line lengths assumed in the analysis, which were vetted by the participating TSO.

Table 1. Transmission line distances between the substations.

<table>
<thead>
<tr>
<th>Transmission Line</th>
<th>Primorska</th>
<th>CERS Moldova</th>
<th>Vulcanesti</th>
<th>Isaccea</th>
</tr>
</thead>
<tbody>
<tr>
<td>Distance (km)</td>
<td>50</td>
<td>200</td>
<td>175</td>
<td>60</td>
</tr>
</tbody>
</table>

The possible connection points (i.e., candidates) in Moldova are: 400 kV CERS Moldova and Vulcanesti Substations. Summary of the substations along the RUM transmission with the corresponding voltage levels are given in Table 2.

Table 2. Summary of the substations along the RUM transmission corridor.

<table>
<thead>
<tr>
<th>Country</th>
<th>Substation</th>
<th>Abbreviation</th>
<th>Existing DC Voltage Level in Substations</th>
</tr>
</thead>
<tbody>
<tr>
<td>Ukraine</td>
<td>Primorska</td>
<td>UKR</td>
<td>✓ 750kV, x 400kV, ✓ 330kV</td>
</tr>
<tr>
<td>Moldova</td>
<td>CERS Moldova</td>
<td>MDV_1</td>
<td>✓ 750kV, ✓ 400kV, ✓ 330kV</td>
</tr>
<tr>
<td>Moldova</td>
<td>Vulcanesti</td>
<td>MDV_2</td>
<td>x 750kV, ✓ 400kV, ✓ 330kV</td>
</tr>
<tr>
<td>Romania</td>
<td>Isaccea</td>
<td>ROM</td>
<td>✓ 750kV, x 400kV, x 330kV</td>
</tr>
</tbody>
</table>

Considering the candidate substations in Moldova and the available voltage levels, three groups of variants are generated as alternatives for the investigations, as presented in Table 3.
Table 3. Substation and voltage level variants to be investigated.

<table>
<thead>
<tr>
<th>Seasonal Variants</th>
<th>Substation Variants</th>
<th>Voltage Level Variants</th>
<th>Connection Type Variants</th>
</tr>
</thead>
<tbody>
<tr>
<td>Summer Peak Load</td>
<td>UCR - MV1, 1 - ROM</td>
<td>2 x 400 kV</td>
<td>400 kV</td>
</tr>
<tr>
<td>Winter Peak Load</td>
<td>UCR - MV2, 2 - ROM</td>
<td>1 x 400 kV</td>
<td>230 kV</td>
</tr>
</tbody>
</table>

Initial assumption, was to analyze a total of 36 scenarios (Substation Variants (4) x Voltage Level Variants (3) x Seasonal Variants (3) = 36) as given in Table 3. However, the initial analysis indicated a strong dependency of results to “Connection Type Variants”. Hence “Connection Type Variants” are also included in the analysis creating a total of 108 scenarios (108 = 36 x 3 (Connection Type Variants)) to be analyzed.

II. METHODOLOGY

The approach in load flow and N-1 contingency analysis will be to search for the maximum amount of power that can be transferred safely from/to Ukraine-Moldova to/from Romania, for each combination of Substation and Voltage Level Variants shown in Table 3. Flowchart of this approach (i.e., algorithm) is given in Fig. 3. Nevertheless, the OPF and contingency analysis are performed and results are recorded for both considering and ignoring the voltage constraints. The effect of voltage constraints on total generation cost and optimum power exchange amounts in Scenario 1 is illustrated in Fig. 4, as an example.

Fig. 4. Total generation cost results of Scenario 1 (with and without voltage constraints).

As illustrated in Fig. 4, the OPF algorithm forces the solution to a higher cost in order to be able to satisfy the voltage constraints. The effect is more observable as the exchange approaches higher values at both directions. However, as the voltage constraints are relaxed, the cost reduces as OPF does not consider voltage constraints. The algorithm of the methodology in Fig. 3 starts from an N-secure case with zero exchanges and iteratively increases the exchange power through the RUM corridor in order to find the optimum power exchange between the parties. The algorithm repeats the following iteratively:

- Assigns an exchange from/to Ukraine-Moldova to/from Romania as a constraint to OPF,
- Performs an OPF to determine the dispatching,
- Compare total generation cost with zero power exchange case in order to determine the realistic transaction limit due to price difference,
- Creates a load flow scenario based on OPF solution,
- Performs N-1 contingency analysis,
- Records the N-1 security violations of each scenario, if any,
- Power exchange is increased in 50 MW steps.
- The analysis performed with and without voltage constraints and the results are compared.

Trading scenarios between RUM countries, as predicted by each party, are given in Table 5. The trading amounts presented in this table are utilized as “indicative” parameters in the analysis. As described above, the approach in OPF analysis is to search for the “maximum” amount of power that can be transferred N-securely between the RUM countries for each scenario that are shown in Table 3. In other words, the algorithm given in Fig. 3 will give the upper limit (i.e., maximum) for the N-secure power trading among RUM countries. The upper limit for the N-secure power trading could be less or more than the corresponding indicative power transfer amounts that are shown in Table 5 (last column).

Table 5. Trading scenarios initially predicted by RUM parties.

<table>
<thead>
<tr>
<th>Country</th>
<th>750 kV</th>
<th>500 kV</th>
<th>400 kV</th>
<th>230 kV</th>
<th>130 kV</th>
</tr>
</thead>
<tbody>
<tr>
<td>Romania</td>
<td>+/-5%</td>
<td>+/-5%</td>
<td>+/-5%</td>
<td>+/-5%</td>
<td>+/-5%</td>
</tr>
<tr>
<td>Ukraine</td>
<td>+/-5%</td>
<td>+/-5%</td>
<td>+/-5%</td>
<td>+/-5%</td>
<td>+/-5%</td>
</tr>
<tr>
<td>Moldova</td>
<td>+/-5%</td>
<td>+/-5%</td>
<td>+/-5%</td>
<td>+/-5%</td>
<td>+/-5%</td>
</tr>
</tbody>
</table>
III. COST/BENEFIT ANALYSES

The OPF analyses cover the largest part of the analysis and create a basis for the cost/benefit analysis which is described in this section. In this section, performance indicators for economic and financial analysis, determination of necessary investments for the corresponding investment scenarios, and calculation of per unit investment and operation and maintenance costs are described.

III.1. PERFORMANCE INDICATORS (IRR, NPV, AND B/C RATIO)

The economic and financial analysis is based on the results of the OPF analysis, which calculates the total savings to region at the optimum power exchange in each loading hour (i.e., winter max, summer max, and summer min). The Cost/benefit analysis was made by comparing the results of the OPF analysis with the investment cost (Inv cost) and operational and maintenance costs (O&M cost) of the candidate investments.

For each scenario, annual cash flow tables for 30 years were determined to conduct the following analyses:

- Internal rate of return (IRR) analysis,
- Net present value (NPV) analysis,
- Benefit/Cost ratio (B/C ratio) analysis.

A 30 year useful life of equipment was assumed for the purposes of the economic/financial calculations. The following parameters (KEPs) were utilized:

- Interest rate of borrowing money for total investment cost
- Loan period
- Discount rate for calculating NPVs

The costs of each investment scenario includes:

- Total investment cost (TIC) at the initial year (USD),
- Annual O&M cost (USD/year).

As described below, annualized savings are considered in the economic and financial analysis. As described in Section 9, sensitivity analyses were performed for the key economic parameters including both AWFs and different generation levels of wind power plants in Romania.

The year 2012 is considered to be the base year in unit costs of the equipment. The annual cash flow table is provided in Table 6.

It should be noted that annual savings can be negative in some investment scenarios that correspond to the most constrained loading conditions (e.g., winter max), even during zero exchange among the countries. This can occur when voltage constraints combine with high technical losses, and higher generation levels of the most costly power plants in the RUM countries than might be dispatched if there was no interconnection between the RUM countries. The total saving is assumed to be zero in such cases because the RUM transmission corridor circuit can be opened to curtail electricity flow in such circumstances.

Table 6. Annual cash flow table of each scenario.

<table>
<thead>
<tr>
<th>Scenarios</th>
<th>Import from</th>
<th>Export to</th>
<th>Transfer</th>
</tr>
</thead>
<tbody>
<tr>
<td>1</td>
<td>Ukraine</td>
<td>Romania</td>
<td>300 MW</td>
</tr>
<tr>
<td>2</td>
<td>Ukraine</td>
<td>Romania</td>
<td>500 MW</td>
</tr>
<tr>
<td>3</td>
<td>Moldova</td>
<td>Romania</td>
<td>480 MW</td>
</tr>
<tr>
<td>4</td>
<td>Ukraine</td>
<td>Romania</td>
<td>1500 MW</td>
</tr>
<tr>
<td>5</td>
<td>Moldova</td>
<td>Romania</td>
<td>480 MW</td>
</tr>
<tr>
<td>6</td>
<td>Ukraine</td>
<td>Romania</td>
<td>1500 MW</td>
</tr>
<tr>
<td>7</td>
<td>Ukraine</td>
<td>Romania</td>
<td>1500 MW</td>
</tr>
<tr>
<td>8</td>
<td>Ukraine</td>
<td>Romania</td>
<td>1500 MW</td>
</tr>
<tr>
<td>9</td>
<td>Romania</td>
<td>Romania</td>
<td>100 MW</td>
</tr>
<tr>
<td>10</td>
<td>Romania</td>
<td>Romania</td>
<td>1500 MW</td>
</tr>
<tr>
<td>11</td>
<td>Romania</td>
<td>Romania</td>
<td>1500 MW</td>
</tr>
</tbody>
</table>

In the LF and OPF analysis, the maximum savings in each scenario are determined at three loading hours along the year (i.e., winter max, summer max, and summer min loading hours). The savings at these loading hours are utilized in determining the annual average savings (i.e., annualization of the savings). Annualization of the savings is based on the annualized weighting factors (AWF) of these three loading hours.

The regional system coincident annual hourly load recordings for 2010 were utilized to determine the AWFs. Annual hourly coincident regional load and its distribution along one year are given in Fig. 5 and Fig. 6, respectively. The following loading hours are indicated in the figures:

- System peak (i.e., winter max)
- System off-peak (i.e., summer min)
- System peak during off-peak season (i.e., summer max)

The AWFs of these three loading hours are indicated in Fig. 6 and summarized in Table 7. Note that, total energy consumed along the year (the area below the blue curve in Fig. 5) is equal to (1):

$$4.29\% \times \text{System peak} + 27.58\% \times \text{System off-peak} + 68.13\% \times \text{System peak during off-peak season}$$

Table 7. AWFs assumed for RUM countries.
This approach is analyzed below for the following parameter and investment scenario (Case_VC-I_W30%):

- **Investment Scenario No: 1**
  - 1x400kV Ukraine-MDV_2-ROM (connection through HVDC B2B substation at Romania)
  - Wind generation level at Romania: 30%
  - Voltage constraints: Ignored
  - Loading Scenarios:
    - Scenario 4: System peak (Winter max)
    - Scenario 3: System off-peak (Summer min)
    - Scenario 2: System peak during off-peak season (Summer max)

The savings which are determined by OPF analyses for the three loading scenarios are given in Table 8.

### Table 8. Annualization of savings for the Scenarios 2, 3 and 4

As illustrated in the table:
- The maximum saving occurs at “system peak” (i.e., winter max).
- The room for OPF is maximum given high generation levels of cost-ineffective power plants in the region.
- The minimum saving occurs at “system off-peak” (i.e., summer min).
- The potential for optimization is minimal due to system constraints at minimum loading conditions
- The availability of cost effective generator capacity in the system is minimum.
- In order to determine the annualized total saving, availability of the line should be estimated (downtime for maintenance and unavailability of the line due to faults must be estimated). An availability of 8322 hours, which corresponds to ≈95% of the hours in a year, is assumed for the economic/financial analysis.
- Annual saving for this investment scenario is calculated as 43,596,276 USD/year, as illustrated in Table 8.

This approach was employed in for investment scenarios to determine the annualized total saving.

### III.2. NECESSARY INVESTMENTS AND CORRESPONDING COSTS

This section reviews the approach to determining the total amount of equipment that should be installed to support each investment scenario. The following assumptions are made:

- Each substation was equipped with a spare bay at the corresponding voltage level, in case of emergency/maintenance/etc.: 750/400 kV or 750/330 kV transformers at the corresponding substations to satisfy n-1 reliability criteria;
- In the scenarios in which there are two substations in Moldova, there will be only one transformer in each substation. This meant that n-1 contingency was satisfied by the transformer in the other substation; and
- As the intermediary substations in the corridor, the new substations in Moldova were assumed to have additional bays– the total number of which depends on the connection type.

The determination of the necessary equipment for different investment scenarios is described in the following subsections.

#### III.2.1. UKRAINE – ROMANIA (1X750 kV AC)

In this investment scenario, the following investments are assumed in Ukraine and Romania:

- **Primorska/Ukraine:**
  - Two 750/330 kV transformers (1250 MVA)
  - Four 750 kV bay:
    - One for the transmission line;
    - Two for the connection of transformers; and
    - One for spare.
  - Three 330 kV bay:
    - Two for the connection of transformers; and
    - One for spare.

- **Isaccea/ Romania:**
  - Two 750/400 kV transformers (1250 MVA).
  - Four 750 kV bay:
    - One for the transmission line;
    - Two for the connection of transformers; and
    - One for spare.
  - Three 400 kV bay:
    - Two for the connection of transformers; and
    - One for spare.

The necessary equipment is summarized in Table 9 below.

### Table 9. Total amount of equipment necessary for Scenario I.
- One for spare.
  - Three 330 kV bay:
    - Two for the connection of transformers; and
    - One for spare.
- CERS Moldova/Moldova:
  - Two 750/400 kV transformers (1250 MVA).
  - Five 750 kV bay:
    - One for the transmission line input;
    - One for the transmission line output;
    - Two for the connection of transformers; and
    - One for spare.
  - Three 400 kV bay:
    - Two for the connection of transformers; and
    - One for spare.
- Isaccea/Romania:
  - Two 750/400 kV transformers (1250 MVA).
  - Four 750 kV bay:
    - One for the transmission line;
    - Two for the connection of transformers; and
    - One for spare.
  - Three 400 kV bay:
    - Two for the connection of transformers; and
    - One for spare.

The necessary equipment is summarized in Table 10 below.

### Table 10. Total amount of equipment necessary for Scenario 2.

<table>
<thead>
<tr>
<th>Scenario</th>
<th>Transmission Corridor</th>
<th>Voltage Level</th>
</tr>
</thead>
<tbody>
<tr>
<td>2</td>
<td>LUK - MDV - ROM.</td>
<td>1 x 750 kV</td>
</tr>
<tr>
<td>Primorska</td>
<td>750 kV bay for transmission line to ROM</td>
<td>1</td>
</tr>
<tr>
<td></td>
<td>750 kV spare bay for RUM</td>
<td>1</td>
</tr>
<tr>
<td></td>
<td>750 kV bay for 750/330 kV power transformer</td>
<td>2</td>
</tr>
<tr>
<td></td>
<td>330 kV bay for power transformers</td>
<td>2</td>
</tr>
<tr>
<td></td>
<td>330 kV spare bay</td>
<td>1</td>
</tr>
<tr>
<td></td>
<td>Total 750 kV bay</td>
<td>4</td>
</tr>
<tr>
<td></td>
<td>Total 330 kV bay</td>
<td>3</td>
</tr>
<tr>
<td></td>
<td>Total 750/330 kV power transformers</td>
<td>2</td>
</tr>
<tr>
<td></td>
<td>750 kV bay for transmission line to ROM</td>
<td>2</td>
</tr>
<tr>
<td></td>
<td>750 kV spare bay for RUM</td>
<td>1</td>
</tr>
<tr>
<td></td>
<td>750 kV bay for 750/400 kV power transformer</td>
<td>2</td>
</tr>
<tr>
<td></td>
<td>400 kV bay for power transformers</td>
<td>2</td>
</tr>
<tr>
<td></td>
<td>400 kV spare bay</td>
<td>1</td>
</tr>
<tr>
<td></td>
<td>Total 750 kV bay</td>
<td>5</td>
</tr>
<tr>
<td></td>
<td>Total 400 kV bay</td>
<td>3</td>
</tr>
<tr>
<td></td>
<td>Total 750/400 kV power transformers</td>
<td>2</td>
</tr>
<tr>
<td>Isaccea</td>
<td>750 kV bay for transmission line to ROM</td>
<td>1</td>
</tr>
<tr>
<td></td>
<td>750 kV spare bay for RUM</td>
<td>1</td>
</tr>
<tr>
<td></td>
<td>750 kV bay for 750/400 kV power transformer</td>
<td>2</td>
</tr>
<tr>
<td></td>
<td>400 kV bay for power transformers</td>
<td>2</td>
</tr>
<tr>
<td></td>
<td>400 kV spare bay</td>
<td>1</td>
</tr>
<tr>
<td></td>
<td>Total 750 kV bay</td>
<td>4</td>
</tr>
<tr>
<td></td>
<td>Total 400 kV bay</td>
<td>3</td>
</tr>
<tr>
<td></td>
<td>Total 750/400 kV power transformers</td>
<td>2</td>
</tr>
</tbody>
</table>

III.2.3. UKRAINE – MOLDOVA _1_ – MOLDOVA _2_ - ROMANIA (1x750 kV)

- Primorska/Ukraine:
  - Two 750/330 kV transformers (1250 MVA).
  - Four 750 kV bay:
    - One for the transmission line;
    - Two for the connection of transformers; and

The necessary equipment is summarized in Table 11 below.

### Table 11. Total amount of equipment necessary for Scenario 4.

<table>
<thead>
<tr>
<th>Scenario</th>
<th>Transmission Corridor</th>
<th>Voltage Level</th>
</tr>
</thead>
<tbody>
<tr>
<td>4</td>
<td>LUK - MDV - ROM.</td>
<td>1 x 750 kV</td>
</tr>
<tr>
<td>Primorska</td>
<td>750 kV bay for transmission line to ROM</td>
<td>1</td>
</tr>
<tr>
<td></td>
<td>750 kV spare bay for RUM</td>
<td>1</td>
</tr>
<tr>
<td></td>
<td>750 kV bay for 750/330 kV power transformer</td>
<td>2</td>
</tr>
<tr>
<td></td>
<td>330 kV bay for power transformers</td>
<td>2</td>
</tr>
<tr>
<td></td>
<td>330 kV spare bay</td>
<td>1</td>
</tr>
<tr>
<td></td>
<td>Total 750 kV bay</td>
<td>4</td>
</tr>
<tr>
<td></td>
<td>Total 330 kV bay</td>
<td>4</td>
</tr>
<tr>
<td></td>
<td>Total 750/330 kV power transformers</td>
<td>2</td>
</tr>
<tr>
<td></td>
<td>750 kV bay for transmission line to ROM</td>
<td>2</td>
</tr>
<tr>
<td></td>
<td>750 kV spare bay for RUM</td>
<td>1</td>
</tr>
<tr>
<td></td>
<td>750 kV bay for 750/400 kV power transformer</td>
<td>2</td>
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<td></td>
<td>400 kV bay for power transformers</td>
<td>2</td>
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<tr>
<td></td>
<td>400 kV spare bay</td>
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<tr>
<td></td>
<td>Total 750 kV bay</td>
<td>5</td>
</tr>
<tr>
<td></td>
<td>Total 400 kV bay</td>
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</tr>
<tr>
<td>CERS Moldova</td>
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<tr>
<td></td>
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<tr>
<td></td>
<td>750 kV bay for 750/400 kV power transformer</td>
<td>1</td>
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<td></td>
<td>400 kV bay for power transformers</td>
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<tr>
<td></td>
<td>400 kV spare bay</td>
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<tr>
<td></td>
<td>Total 750 kV bay</td>
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<tr>
<td></td>
<td>Total 400 kV bay</td>
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<tr>
<td></td>
<td>Total 750/400 kV power transformers</td>
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</tr>
<tr>
<td>Vulcanești</td>
<td>750 kV bay for transmission line to ROM</td>
<td>1</td>
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<tr>
<td></td>
<td>750 kV spare bay for RUM</td>
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<tr>
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<td>750 kV bay for 750/400 kV power transformer</td>
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<tr>
<td></td>
<td>400 kV spare bay</td>
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<tr>
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<td>Isaccea</td>
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<td>Total 750 kV bay</td>
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<tr>
<td></td>
<td>Total 750/400 kV power transformers</td>
<td>2</td>
</tr>
</tbody>
</table>
IV. HVDC BACK TO BACK CONNECTION ANALYSES

In this section, first the challenges with the HVAC interconnection of ENTSO-E and IPS/UPS systems are discussed. Then, the analysis for the HVDC interconnections of RUM Countries is presented. Romania is connected to the ENTSO-E system, whereas Ukraine and Moldova are connected with IPS/UPS network as shown in Fig.6.

![ENTSO-E and IPS/UPS systems at the RUM countries’ area](image)

Following the EU-Russia energy dialogue, an extensive study was launched in 2004 under UCTE guidance with the aim of identifying the technical and operational preconditions for the interconnection of the two largest European power systems – UCTE (now ENTSO-E) and IPS/UPS. The possibility of the interconnection of the two systems, which would allow for direct technical and commercial cooperation in the field of electric power, was investigated. Particular attention was devoted to stability of interconnected networks, the prevention of crisis situations and the legal aspects of such interconnected operation. The results of the study proved the technical possibility of such interconnection, but concluded that the investment required to operate the system in a secure and stable manner were prohibitive.

Given the cost and technical challenges associated with AC connection of the ENTSO-E and IPS/UPS networks, HVDC back-to-back technology interconnection of the RUM countries was considered as a variant for the short/mid-terms.

V. HVDC BACK TO BACK TECHNOLOGIES

There are two primary HVDC Back to Back technologies: Line Commutated Converter (LCC) and Voltage Source Converter (VSC). HVDC Back to Back substations based on conventional Line Commutated Converter (LCC) technology depend on the Short Circuit MVA (SCMVA) at the connection point to the grid. The new VSC technology substations can operate independent from the SCMVA at the connection point. Today, both technologies are being deployed.

The chronological development of the two HVDC technologies is given in Fig.7.

![Chronological development in LCC and VSC](image)

V.1. REQUIREMENT OF HARMONIC FILTERS

LCC based HVDC Back to Back substations generally require harmonic filters with a capacity of almost 60% of the substation [2]. For example, for a 300 MW block substation, the capacity of the necessary harmonic filters is 300*0.6 = 180 MVar.

V.2 ESCR CRITERIA

The results of the ESCR calculation results are presented in this section to determine the acceptable level of the LCC technology based HVDC Back to Back substations.

\[
ESCR = \frac{SCMVA_{grid} - S_{filter}}{P_{dc}}
\]

(2)

In this formula, the contribution of the filters to SCMVA is subtracted to consider the true SCMVA of the grid. In this article, HVDC Back to Back connection is modelled by splitting the networks at the point of HVDC connection and introducing POSITIVE and NEGATIVE loads at appropriate sides. The schematic representation of such modelling is illustrated in Fig.7. As seen in the figure, the power flow through HVDC Back to Back substation from Primorska to Isaccea is modelled by splitting the networks and introducing a POSITIVE Load at Primorska side and a NEGATIVE Load at Isaccea side. It should be noted that a NEGATIVE Load is preferred in representing power injection rather than modelling a generator, in order to avoid unrealistic reactive support from

<table>
<thead>
<tr>
<th>LCC Technology</th>
<th>VSC Technology</th>
</tr>
</thead>
<tbody>
<tr>
<td>0.08 Euro/VA</td>
<td>0.11 Euro/VA</td>
</tr>
</tbody>
</table>

![Line Commutated Converter (LCC)](image)

![Voltage Source Converter (VSC)](image)
the HVDC Back to Back via the generator. Given this representation, the SCMVA contribution of the HVDC Back to Back filters is not considered in the load flow and short circuit analysis. Therefore, the ESCR should be calculated as in (3).

$$ESCR = \frac{SCMVA_{grid}}{P_{dc}}$$  

(3)

For the secure operation of HVDC Back to Back substation that is based on LCC the

$$ESCR \geq 3 \text{ (base case)}$$  

(4)

Essentially, the ESCR is different at each connection point of the HVDC Back to Back substations given different topologies. For the sake of security, the minimum value among the SCMVA at each connection point is considered in calculating of the ESCRs. The available HVDC Back to Back substation capacity is calculated assuming that total capacity of the substation is formed by 300 MW blocks, while taking into account the ESCR criteria (4).

V.3 DETERMINATION OF TOTAL CAPACITY OF HVDC BACK TO BACK SUBSTATION

It is assumed that the HVDC Back to Back substation blocks will be in the order of 300 MW capacities. The following arguments support this approach:

- 300 MW capacity HVDC Back to Back substations are available in the market.
- The order of 300 MW is plausible to match the optimum substation capacity with the optimum power exchange amounts that are determined in LF (Load Flow) and OPF (Optimal Power Flow) analysis.

For example, the approach in determining the total capacity of the HVDC Back to Back substation is presented below (1x400 kV transmission line between Ukraine - Romania through HVDC Back to Back substation in Ukraine):

- Loading condition of the scenario: Summer maximum.
- Wind generation level in Romania: Normal (i.e., generation level of the wind power plants in Dubrudja-ROM region is 30% of the capacity).
- OPF results at base case (i.e., ignoring N-1 contingency):
  - 700 MW (Ukraine => Romania)
  - N-1 security exchange technical limit:
    - 1,300 MW (Ukraine => Romania)
    - Since 700 < 1300, 700 MW power exchange is feasible in the sense of N-1 security concern.
- Voltage collapse power exchange limit:
  - 1,500 MW (Ukraine => Romania)
  - Since 700 < 1,500, 700 MW power exchange is feasible in the sense of voltage collapse concern.
- Assuming that HVDC Back to Back substation is composed of 300 MW blocks, total number of block to realize 700 MW power exchange is three (3*300 = 900 > 700)
  - Total capacity of the HVDC BACK TO BACK substation is 900 MW.
- ESCR criteria:
  - Maximum SCMVA of the grid at the HVDC Back to Back substation is calculated as 2.063 MVA
  - ESCR = 2.063/900 = 2.29
  - Since 2.29 < 3, total capacity of 900 MW is NOT acceptable in the sense of ESCR criteria.
  - If one block among the three blocks is removed, then the total capacity of the substation is 2x300 = 600 MW

$$ESCR = \frac{2,063}{600} = 3.43 > 3 \Rightarrow \text{acceptable}$$

An HVDC Back to Back capacity of 600 MW is proven in the summer maximum loading conditions. Similar analyses were performed for winter maximum and summer minimum loading conditions, as well. The total capacity of the HVDC Back to Back substation is considered to be the maximum capacity determined among three loading scenarios. This approach is considered in all scenarios that include HVDC Back to Back substation.

VI. CONCLUSIONS

Voltage constraints were local problems that could be resolved through network operations in the short term (e.g., proper selection of generator voltage set points) and relatively inexpensive capacitor/reactor investments in the mid-term. Hence, voltage constraints are ignored in certain cases to determine the maximum volume of power exchange among the countries. The maximum voltage deviation at the key nodes was observed to be +/-20%, which could be resolved by proper compensation through the provision of additional reactors.

The increase in wind generation in Romania dramatically limited the ability of the RUM countries to optimize the regional generation fleet based on the cost of production. In some investment scenarios, the flow of power changed direction from north → south to south → north when the wind power plant generation in Romania increased from 30% to 70% and it is designated as must run. This occurs when the OPF algorithm forced inefficient high cost generators, first in Romania and then in Moldova and Ukraine, to reduce their generation in favor of must run wind. This process continued until the reduction of generation in Ukraine and Moldova became so much more cost effective than the reduction of generation in Romania that the power flow changed direction. From this point onward, Romania began exporting power in a northward direction to Moldova and Ukraine.

It is important to note that for the investment scenario of a 400 kV connection passing through a HVDC B2B substation, the benefit/cost ratio was > 1, when Romanian must run wind generation was modeled with a 30% capacity factor.

Connection through the HVDC Back to Back was superior to connection through AC options in almost every investment scenario considered. This was because the HVDC connection reduced technical network constraints to increase power exchange, enlarging the scope for power flows in the sub-region.

In fact, HVDC B2B was the only investment solution which resulted in benefit/cost > 1 when considering the scenario of Romanian must run operating with a 30% capacity factor. And, the technical challenges to synchronizing the current IPS/UPS and ENTSO-E members of the RUM working group would inhibit interconnection via high voltage AC interconnections for the foreseeable future. Therefore, HVDC technology based interconnection of the RUM countries seemed the most rational solution in the short/mid-term.

There was no significant difference revealed in the cost/benefit analyses for the different investment scenarios related to the configuration of the corridor, i.e., either directly from Ukraine to Romania or through Moldova. If the interconnection between RUM countries were realized in intermediate steps, (for example, if the connection between Romania and Moldova were realized before all three countries are interconnected), energy trade between Romania and Moldova could begin before the trading among all three countries by directing a generator in Moldova to operate synchronously with Romania in island mode.
REFERENCES


[3]. “Feasibility study of the asynchronous interconnection between the power systems of Iran and Turkey”, Interim Report, Professor Francesco Iliceto, Consultant to TEIAS, Rome (Italy), November 2009.

* Siemens PTI Documentation
Abstract—To deal with disease, one should distinguish between the normal value and abnormal values related to disease manifestations. Beside sign and symptom, biomarkers are among the milestone tools in human disease diagnosis and treatment including strategies for further drug development. Nowadays, there is significant increase in quantitity and quality of biomarkers used in the process of diagnosis/differential diagnosis, prognosis, treatment decision making and therapy outcomes monitoring. This symposium is aimed to highlight advances and limitations of biomarkers use to date. Biomarkers have an increasingly important clinical role in managing patients with heart failure as well as those with kidney disease, both common conditions with generally poor prognostic outcomes and huge impacts on healthcare economics. For patients with chronic heart failure, biomarkers have become centre place in streamlining diagnostic pathways as well as identifying those with worse prognosis. There is much interest in the role for biomarkers in identifying patients at risk of acute kidney injury, although a number of these currently remain as research tools or are in the early stages of evaluation in clinical practice. Patients with cardiorenal syndrome represent a particular challenge to the clinician, and recent studies have suggested a valuable clinical role for certain biomarkers in this setting, either on their own or in combination. This paper is aimed to highlight advances and limitations of biomarkers use to date focusing on biomarkers with a current clinical role in patients with cardiorenal disease (natriuretic peptides and neutrophil gelatinase-associated lipocalin), although brief reference will be made to other biomarkers with potential future application including those paralleling the drug toxicity.

Keywords—Biomarkers, Diagnosis, Prognosis, Treatment Guide

I. INTRODUCTION

A biomarker, or biological marker, generally refers to a measurable characteristic, which may be used as an indicator of certain biological state or condition. The term occasionally also refers to a substance whose presence indicates the existence or absence of some signs and symptoms in living organisms or a sample drawn from. Human physiology is complex and multifactorial and exhibits the properties of a system; where the endocrine system manages metabolism, which is the basis of the continuity of life and the metabolic activity managed by the endocrine system results in the output of biomarkers that reflect the functional achievement of specific aspects of metabolism; while biomarkers are related to each other in ratios, it contextualizes one type of function relative to another to which it is linked anatomically, sequentially, chronologically, biochemically, etc. [1] Thus, biomarkers are used in many scientific fields. Biomarkers in medicine are often measured and evaluated to examine normal biological processes, pathogenic processes, or pharmacologic responses to a therapeutic intervention. In medicine, a biomarker can also be a traceable substance that is introduced into an organism as a means to examine organ function or other aspects of health. For example, rubidium chloride is used as a radioactive isotope to evaluate perfusion of heart muscle. It can also be a substance whose detection indicates a particular disease state, for example, the presence of an antibody may indicate an infection. Before the development of dentic testing methods for example
only taste of salt of baby meant that the the bayby dies soon at the time care for cystic fibrosis was not so progressive as today. Eventoday the first routine green test for diagnosis of cystic fibrosis is simple sweat-test to estimae the extent of CI- ions loss in the sweat typcal marker for the gentical disease frequent in caucacian race. More specifically, a biomarker indicates a change in expression or state of a protein that correlates with the risk or progression of a disease, or with the susceptibility of the disease to a given treatment. Biochemical biomarkers are often used in clinical trials, where they are derived from bodily fluids that are easily available to the early phase researchers. A useful way of finding genetic causes of diseases such as schizophrenia has been the use of a special kind of biomarker called an endophenotype. Other biomarkers can be based on measures of the electrical activity of the brain (EEG) or volumetric measures of certain brain regions (using magnetic resonance imaging) or saliva testing of natural metabolites, such as saliva nitrite, a surrogate marker for nitric oxide. Biomarker is generally a molecule that is measured as a marker of normal biological processes, disease processes or the response to a diagnostic or therapeutic intervention. Kidney diseases like acute kidney injury, chronic kidney disease, diabetic nephropathy, glomerular disease, renal cancer and preeclampsia still have a high morbidity. Measurement of biomarkers in the blood or urine that detect patients at risk of kidney diseases or that detect kidney diseases in the earliest stage may ultimately result in preventative or earlier or more effective treatments for kidney diseases. The use of the term "biomarker" has been dated back to as early as 1980. [2] In relevance with its application a biomarker is defined as a characteristic that is objectively measured and evaluated as an indicator of normal biological processes, pathogenic processes, or pharmacologic responses to a therapeutic intervention.[3]

In general, the use of biomarkers to aid diagnosis and treatment is increasing rapidly as genomics and proteomics help us expand the number of markers we can use and as an improved understanding of the pathophysiology of diseases guides their use. However, as with all rapidly expanding fields, there is the risk of excessive enthusiasm unless we are circumspect about the data that guide the clinical use of these new tools. This paper focuses first on how to use several biomarkers, which at present are the best validated of the new markers, and will hopefully provide insight into how to use this biomarker more productively by distinguishing subsets of patients and by providing an understanding of the meaning of elevations in various clinical situations individually.

![Fig. 1. An ideal biomarker should incorporate the characteristics illustrated in the figure.](image)

II. RENAL MARKERS

Over the past few years and with the use of innovative genomic and proteomic tools, several molecules that their urinary concentration is modified during acute kidney injury have been identified and proposed as biomarkers. Among the most studied biomarkers are neutrophil gelatinase-associated lipocalin-2, kidney injury molecule-1, interleukin-18, cystatin C, N-acetyl-β-D-glucosaminidase, liver fatty-acid binding protein, and heat shock protein 72. Recently review on several existing and recently used renal biomarkers and compared the sensitivity and specificity of each biomarker for the appropriate diagnosis of acute kidney injury, as well as its ability to stratify renal injury and to monitor a renoprotective pharmacologic strategy confirming their importance [5] Acute kidney injury (AKI) has a number of triggers, including ischaemia, nephrotoxins, radiocontrast, and bacterial endotoxins. It occurs in around one-third of patients treated in intensive care.
unit (ICU) and is even more prevalent in cardiac surgery patients. There is a higher mortality in patients with AKI compared with non-AKI counterparts, and in severe AKI requiring renal support, the 6 month mortality is >50%. Unlike the progressive development of biomarkers in cardiology, there have been few changes in kidney diagnostic markers. Creatinine is still used as an indicator of kidney function but not of the parenchymal kidney injury. Serum creatinine (sCr) concentration does not change until around 50% of kidney function is lost, and varies with muscle mass, age, sex, medications, and hydration status. The lag time between injury and loss of function, risks missing a therapeutic opportunity, and may explain the high associated mortality. Novel biomarkers of AKI- and failure include neutrophil gelatinase-associated lipocalin, N-acetyl-β-d-glucosaminidase, kidney injury molecule-1, interleukin-18, and cystatin C. The pathophysiology associated with accumulation of these markers in plasma and urine is not clear, but a common denominator is inflammation. Some of these new AKI biomarkers may have clinical applicability in anaesthesia and intensive care in the future. It is possible that a 'kidney biomarker panel' will become standard before and after major surgery. If elevated or positive, the anaesthetist must take special care to optimize the patients after operation on the surgical wards or ICU to avoid further nephrotoxic insults and initiate supplementary care. [6] Acute kidney injury usually refers to the rapid loss of renal function. In clinical practice, AKI is common among hospitalized patients of all age groups including neonates and remains an important cause of morbidity and mortality due to its late diagnosis and therefore delayed therapeutic intervention. Although the precise incidence of AKI in newborns is unknown, several studies have reported that 8 to 24% of all critically ill newborns in neonatal intensive care units may develop the condition. We aim at reviewing the existing literature on novel serum and urinary biomarkers and discuss their role in the early diagnosis and prognosis of AKI in newborns. Specifically, this review will focus on cystatin C (CysC), neutrophil gelatinase-associated lipocalin (NGAL) and interleukin-18 (IL-18) in serum and on CysC, NGAL, kidney injury molecule-1, and IL-18 in urine. [7] Acute graft dysfunction can be caused by ischaemic damage or immunological injury leading to serious consequences both in the short and long term. We are in a desperate need for biomarkers of immune and nonimmune injury at different time points of the transplantation time course, beginning from a potential kidney donors where acute kidney damage can pass unnoticed, during the early post-transplant periods to predict acute transplant dysfunction due to various causes and during long term follow up to predict chronic histological changes. The implementation of these novel biomarkers could increase the sensitivity of diagnosis and monitoring of kidney injury in kidney transplant recipients. Traditionally acute graft dysfunction is diagnosed by measuring serum creatinine concentrations. Unfortunately rise in serum creatinine is a late sign of kidney damage. It indicates rather predicts the damage. The treatment, in order to be effective, must be instituted very early after the initiating insult, well before the serum creatinine even begins to rise. Fortunately, emerging technologies such as functional genomics and proteomics have uncovered novel candidates that are emerging as potentially useful biomarkers of acute kidney injury (AKI). The most promising of biomarkers in AKI for clinical use include a plasma panel consisting of Neutrophil Gelatinase-Associated Lipocalin (NGAL) and Cystatin C and a urine panel including NGAL, IL-18 and Kidney Injury Molecule 1 (KIM-1). Most of these biomarkers were developed in non-transplant AKI, yet their role in clinical transplantation has to be identified. [8] Although heart disease is distinguished organ impairment, it is becoming increasingly recognized that manifestations of congenital heart disease (CHD) extend beyond the cardiovascular system. The factors contributing to renal dysfunction in patients with CHD are multifactorial, with acute kidney injury at time of cardiac surgery playing a major role. Acute kidney injury is often diagnosed based on changes in serum creatinine and estimated glomerular filtration rate (eGFR). Such measurements are often late and imprecise. Recent data indicate that urinary biomarkers interleukin-18 (IL-18) and neutrophil gelatinase-associated lipocalin (NGAL) are earlier markers of AKI. We sought to determine the efficacy of urinary IL-18 and NGAL for detecting early acute kidney injury in patients undergoing surgical pulmonary valve replacement, where both NGAL and IL-18 are
early predictive biomarkers of acute kidney injury, and both increase in tandem after surgical PVR. Importantly, both rise before an increase in creatinine or a decrease in eGFR is present and it has been concluded that monitoring both biomarkers may allow for earlier detection and subsequent interventions to prevent AKI at time of surgery for CHD.[9] Neutrophil gelatinase-associated lipocalin (NGAL) appears to be a promising biomarker for the early diagnosis of acute kidney injury (AKI); however, a wide range in its predictive value has been reported. Meta-analysis of diagnostic test studies using custom-made standardized data sheets recording plasma/serum and urine NGAL within 6 hours from the time of insult (if known) or 24-48 hours before the diagnosis of acute kidney injury if the time of insult was not known were analysed. The primary outcome was acute kidney injury, defined as an increase in serum creatinine level > 50% from baseline within 7 days or contrast-induced nephropathy (creatinine increase > 25% or concentration > 0.5 mg/dL in adults or > 50% increase in children within 48 hours). Other outcomes predicted using NGAL were renal replacement therapy initiation and in-hospital mortality, and the results using a hierarchical bivariate generalized linear model to calculate the diagnostic odds ratio (DOR) and sample size-weighted area under the curve for the receiver-operating characteristic (AUC-ROC), indicated that NGAL level appears to be of diagnostic and prognostic value for acute kidney injury. [10] The early detection of acute kidney injury (AKI) may be become possible by several promising early biomarkers which may facilitate the early detection, differentiation and prognosis prediction of AKI. In this study, we investigated the value of urinary liver-type fatty acid-binding protein (L-FABP), neutrophil gelatinase-associated lipocalin (NGAL) and their combination in predicting the occurrence and the severity of AKI following cardiac surgery. Urinary L-FABP and NGAL increased at an early stage after cardiac surgery indicating that the combination of the two biomarkers enhanced the accuracy of the early detection of postoperative acute kidney injury after cardiac surgery before a rise in SCr. [11]. In some conditions related to exposure to nephrotoxic drugs like aminoglycosides and glycopeptides, antibiotics, of course the most guiding biomarker may be elevation of serum creatinine level.[12], [13] In rare circumstances serum creatinine may paradoxically decline so that one may overlook the underlying renal toxicity hanging on low serum creatinine level, unless the drug level is promptly monitored.[14] Therefore, serum creatinine should carefully evaluated within the context of state of hydration, underlying disease, and nephrotoxic agents exposure in the past or at present.

III. LIVER BIOMARKERS

Liver as the largest metabolically active organ has its own relatively specific biomarkers regarding acute and chronic insults or toxic injury. Serum bilirubin level and several enzymes including aminotransferases are used to assess liver cells integrity. Albumin and other liver products can also mark liver fucinal or synthetic capacity state. The liver markers may be associated also with extrahepatic disorders requiring careful interpretation.

<table>
<thead>
<tr>
<th>Marker</th>
<th>Indicates</th>
<th>Tissue expression</th>
<th>Liver localization</th>
<th>Function</th>
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</thead>
<tbody>
<tr>
<td>Alanine Aminotransferase (ALT)</td>
<td>Injury</td>
<td>Liver, skeletal muscle</td>
<td>Hepatocytes</td>
<td>Glucose-Altaine cycle</td>
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<tr>
<td>Glutamate Dehydrogenase (GLDH)</td>
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<td>Liver, kidney, muscle, intestine</td>
<td>Hepatocytes (cytosol)</td>
<td>Amino acid oxidation, urea production</td>
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<td>Malate Dehydrogenase (MDH)</td>
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<td>Hepatocytes (cytosol)</td>
<td>TCA cycle</td>
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<td>Glutathione S-transferase-α (α-GST)</td>
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<td>Hepatocytes</td>
<td>Glutathione S-transferase</td>
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<tr>
<td>Purine Nucleoside Phosphorylase (PNP)</td>
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<td>Hepatocytes, Endothelial, Kidney</td>
<td>Purine pathway</td>
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<td>Arginase-1 ( Arg-1)</td>
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<td>Hepatocytes</td>
<td>Urea cycle</td>
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<td>Hepatocytes</td>
<td>Enterase protects lipoproteins from lipid peroxidation</td>
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<tr>
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<td>Hepatocytes</td>
<td>a-keto-glyoxalase, dihydroxyacetone phosphatase</td>
</tr>
<tr>
<td>miR122 (maybe others)</td>
<td>Injury</td>
<td>Liver</td>
<td>Hepatocytes</td>
<td>small non-coding RNA expression</td>
</tr>
</tbody>
</table>

Table 1. Short summary of some liver biomarkers and their localization

IV. CURRENTLY USED CARDIAC SPECIFIC BIOMARKERS

Despite better establishment of currently used cardiac-specific serum markers unfortunately non-meets all the criteria for an "ideal" marker of acute myocardial infarction (AMI). No test is both highly
sensitive and highly specific for acute infarction within 6 hours following the onset of chest pain, the timeframe of interest to most emergency physicians in making diagnostic and therapeutic decisions. Patients presenting to the emergency departments (ED) with chest pain or other symptoms suggestive of acute cardiac ischemia therefore cannot make a diagnosis of AMI excluded on the basis of a single cardiac marker value obtained within a few hours after symptom onset. The total CK level is far too insensitive and nonspecific a test to be used to diagnose AMI. It retains its value, however, as a screening test, and serum of patients with abnormal total CK values should undergo a CK-MBmass assay. Elevation in CK-MB is a vital component of ultimate diagnosis of AMI, but levels of this marker are normal in one fourth to one half of patients with AMI at the time of ED presentation. There have been several improvements of CK-MB assay timing and subform quantification that appear highly useful for emergency physicians. Rapid serial CK-MB assessment greatly increases the diagnostic value of the assay in a timeframe suitable for ED purposes but unfortunately still misses about 10% of patients ultimately diagnosed with acute MI. Assays of CK-MB subforms have very high sensitivity, and, although unreliable within 4 hours of symptom onset, have excellent diagnostic value at 6 or more hours after chest pain begins. Automated test assays recently have become available and could prove applicable to ED settings. The cardiac troponins are highly useful as markers of acute coronary syndromes, rather than specifically of AMI, and abnormal values at any time following chest pain onset are highly predictive of an adverse cardiac event. The ED applicability of the troponins is severely limited, however, because values remain normal in most patients with acute cardiac events as long as 6 hours following symptom onset.

Myoglobin appeared promising as a marker of early cardiac ischemia but appears to be only marginally more sensitive than CK-MB assays early after symptom onset and less sensitive than CK-MB at 8 hours or more after chest pain starts. Rapid serial myoglobin assessment, however, appears highly useful as an early marker of AMI. The marker has a very narrow diagnostic window. The clinician is left with several tests that are highly effective in correctly identifying patients with AMI (or at high risk for AMI), but none that can dependably exclude patients with acute coronary syndromes soon after chest pain onset. A prudent strategy when assessing ED patients with chest pain and nondiagnostic ECGs is to order CK-MB and troponin values on presentation in the hope of making an early diagnosis of AMI or unstable coronary syndrome. Although it is recognized that normal values obtained within 6 hours of symptom onset do not exclude an acute coronary syndrome, patients at low clinical risk and having normal cardiac marker tests could be provisionally admitted to low-acuity hospital settings or ED observation. After 6 to 8 hours of symptom duration has elapsed, the cardiac-specific markers are highly effective in diagnosing AMI, and such values obtained can be used more appropriately to make final disposition decisions. At no time should results of serum marker tests outweigh ECG findings or clinical assessment of the patient's risk and stability.[15]

Biochemical markers provide clinicians with an important tool for the assessment of acute coronary syndromes. Biochemical markers, including total creatine kinase (total CK), creatine kinase-MB (CK-MB), the MB isoforms, and myoglobin, as well as the troponins--cardiac troponin T (cTnT) and cardiac troponin I (cTnI)--are all used for assessment of the suspected acute myocardial infarction (AMI) patient. In the context of myocardial infarction (MI) diagnosis, total CK is a relatively sensitive marker, but it lacks myocardial specificity because skeletal muscle contains high concentrations of CK. CK-MB is the benchmark for biochemical markers and has both high sensitivity and specificity; however, CK-MB is also present in skeletal muscle and is not diagnostic until eight to twelve hours after onset of symptoms. The MB isoforms are diagnostic earlier but have the same cardiac specificity issues as CK-MB. Myoglobin becomes abnormal about one hour after onset of symptoms and is a sensitive marker for MI; however, myoglobin is cleared quickly and is not cardiac specific. Furthermore myoglobin together with CK may be significantly elevated in association to some drug overdose leading severe hypokalaemia leading to lifethreatening arrhythmia. [16] Both cTnT and cTnI are cardiac specific and...
show high sensitivity and specificity for MI. Risk stratification of acute coronary syndrome patients is another role for biochemical markers; CK-MB, cTnT and cTnI have all been proposed for this function. Compared with CK-MB, both cTnT and cTnI are better able to predict short-term mortality following the index event. Analysis using a logistic regression model that included the electrocardiogram, cTnT, and cTnI showed that cTnT is the most useful marker for risk stratification and cTnT was reported to be able to predict which patients will benefit from treatment with regimens of low molecular weight heparin. [17] In the past decade studies have suggested that immunoassay of cardiac troponin T (cTnT) provides a more sensitive measurement of myocardial necrosis than creatine kinase MB (CK-MB) mass concentration. Abbas et al. who, compared the release of cTnT and CK-MB isoenzyme in patients undergoing percutaneous coronary angioplasty, and investigated the clinical, procedural, and angiographic correlates of abnormal elevations of both of these markers concluded that > 40% of patients undergoing coronary angioplasty have evidence of minor degrees of myocardial damage, as evidenced by cTnT release. Thus, high-risk coronary lesions and both minor and major complications of angioplasty are associated with cTnT release. cTnT appears to be a more sensitive marker of myocardial injury than CK-MB under these circumstances. In comparison with isolated cTnT rise, elevation of both CK-MB and cTnT may be indicative of greater levels of myocardial injury. [18] Variable degrees of myocardial cell injury during PTCA and stent insertion have been observed, based on rises in creatine kinase MB isoenzyme (CK-MB) and cardiac troponin T (cTnT) 6-24 h post-procedure. As there are many variations in technique within each procedure it would be helpful to be able to determine objectively the degree of myocardial damage in order to optimize technique. Harris et al. who measured CK-MB, cTnT and cardiac troponin I (cTnI) to ascertain which is the most sensitive marker for minor myocardial damage in this setting demonstrated that cTnI was the most sensitive indicator of minor myocardial damage. [19] In the study performed to determine the most sensitive biochemical marker for the detection of cardiac myocyte damage potentially sustained during percutaneous coronary intervention (PCI) and to assess whether such a marker can be used to identify patients at increased risk of poor subsequent clinical outcome, Nageh et al. revealed that cTnI proved to be the most sensitive marker in detecting myocardial necrosis following PCI. CK-MB, cTnT and cTnI all provided similarly reliable prognostic information, with cTnT and cTnI being marginally superior in predicting MACE at follow up. [20] The explosion in cardiovascular biomarkers has been stimulated in large part by proteomics, genomics and an improved understanding of the pathophysiology of cardiovascular disease. In the past few years, major changes have evolved in clinicians’ ability to use cardiac troponin, BNP and CRP which will markedly impact on their clinical utility. The recent changes including greater sensitivity, a better understanding of what specific fragments are being measured and now these markers are used clinically as obligatory, but the search for better markers must continue. [21] Jaffe et al discussed the use of troponin as well as the knowledge gaps associated with emerging biomarkers such as B-type natriuretic peptide and C-reactive protein, which are increasingly moving toward more productive clinical use and reflected on some of the large number of markers that are still in development. [22] The prognostic significance of elevations in creatine kinase-MB and troponin T (cTnT), which have been conventionally measured 6 to 8 h after percutaneous coronary intervention (PCI), has been established. but, the time to peak biomarker appearance in the circulation has not been well defined. Later study revealed that more cTnT than CK-MB elevations occur after PCI; however, both biomarkers demonstrate a longer time to peak value than anticipated in clinical practice. Early surveillance monitoring (< 12 h) does not detect peak biomarker levels, especially in patients with normal baseline values. If peak levels are to be used to determine prognosis, then longer time intervals should be used for post-PCI surveillance. The timing of peak elevations appears to be influenced by baselines values as well. Early elevations may reflect the conjoint effects of injury associated with the disease process and the intervention itself. [23] The adverse prognostic significance of biomarker elevations after percutaneous coronary intervention (PCI) is well established. However, often baseline troponin values are not included in the analysis or sensitive criteria are not employed till Miller et al.
assessed the timing and magnitude of post-PCI troponin T (cTnT) levels and their relationships to outcomes in patients with and without pre-PCI baseline cTnT elevations using a sensitive assay and sensitive cut-off values and finally revealed that long-term prognosis is most often related to the baseline pre-PCI troponin value and not the biomarker response to the PCI.[24]

V. EARLY NEURON DAMAGE BIOMARKERS

Severe brain disorders can be expressed as markedly abnormal encephalopathic EEG patterns in neonates who are usually neurologically depressed, with abnormal levels of reactivity and tone. EEG sleep study can serve as a useful neurophysiologic screening procedure for the child suspected of having a subclinical presentation of an emerging static encephalopathy; longitudinal studies will then document deviations from expected ontogeny in the vulnerable child who is later stressed by environmental and socioeconomic factors. [25] Encephalopathy consequent on perinatal hypoxia-ischemia occurs in a number of newborns and frequently leads to serious and tragic consequences that devastate lives and families, with huge financial burdens for society. Brain pH changes are closely involved in the control of cell death after injury: an alkalosis enhances excitability while a mild acidosis has the opposite effect. Brain alkalosis in babies with neonatal encephalopathy serially studied using phosphorus-31 magnetic resonance spectroscopy studies during the first year after birth (151 studies throughout the year including 56 studies of 50 infants during the first 2 weeks after birth) proved that an alkaline brain pH was associated with severely impaired outcome; the degree of brain alkalosis was related to the severity of brain injury on MRI and brain lactate concentration; and a persistence of an alkaline brain pH was associated with cerebral atrophy on MRI. [26] Although intensive fetal and neonatal resuscitative efforts have reduced the severe expression of the neonatal brain disorder termed hypoxic-ischemic encephalopathy, neonates may alternatively express altered EEG-sleep organization over the first days of life after asphyxia which may mimic mild or moderate hypoxic-ischemic encephalopathy. EEG-sleep studies can assist in a more accurate classification of newborn encephalopathy that does not satisfy the criteria for hypoxic-ischemic encephalopathy. [27] Hypoxic ischemic encephalopathy after perinatal asphyxia is a major cause of mortality and morbidity in infants. Asphyxiated infants displayed neuronal cell damage and reactive glial changes with strong aquaporin-4 immunoreactivity on astroglial cells within hippocampi in 50% of cases, whereas in patients with seizures, the expression of metabotropic glutamate receptors was increased in glial cells, where cases with seizures displayed increased microglial activation and greater expression of the inflammatory markers interleukin 1β and complement 1q compared with those in cases without seizures confirming the complex cascade of cellular and molecular changes occurring in the human neonatal hippocampus after perinatal asphyxia. [28] Anticonvulsant therapy has been used in infants with perinatal asphyxia in order to prevent seizures. However, long term anticonvulsant therapy may lead to inhibition of brain development. Therefore, the routine use of anticonvulsant therapy to prevent seizures following perinatal asphyxia needs to be evaluated. At the present time, anticonvulsant therapy to term infants in the immediate period following perinatal asphyxia cannot be recommended for routine clinical practice, other than in the treatment of prolonged or frequent clinical seizures. Any future studies should be of sufficient size to have the power to detect clinically important reductions in mortality and severe neurodevelopmental disability. [29]

VI. CANCER BIOMARKERS

Among pioneers in the biomarker area of malignant (cancer) diseases, alpha-fetoprotein (AFP) is a well-known diagnostic biomarker used in medicine to detect fetal developmental anomalies such as neural tube defects or Down’s syndrome is also used to date to follow up the development of tumors such as hepatocellular carcinomas todate since its discovery more than half a century ago. [30] Tumor markers in general are endogenous proteins or metabolites whose amounts or modifications are indicative of tumor state, progression characteristics, and response to therapies. They are present in tumor tissues or body fluids and encompass a wide variety of molecules, including...
transcription factors, cell surface receptors, and secreted proteins. Effective tumor markers are in great demand since they have the potential to reduce cancer mortality rates by facilitating diagnosis of cancers at early stages and by helping to individualize treatments. During the last decade, improved understanding of carcinogenesis and tumor progression has revealed a large number of potential tumor markers. It is predicted that even more will be discovered in the near future with the application of current technologies such as tissue microarrays, antibody arrays, and mass spectrometry. To apply these discoveries to patient care, vigorous validation and assay development for many tumor markers is currently underway. The reason behind so few biomarkers reaching the clinic can largely be explained by the inability of current technologies to consistently and quantitatively verify the presence of the candidates in patient samples and the failure, thus far, to identify biomarkers with high specificity for a particular disease as some times none of the cancer biomarkers demonstrate the specificity required for diagnosis when used alone. Therefore, the development of panels of proteins further may be crucial to achieve the specificity required for early cancer diagnosis, and is interesting to speculate that members of such panels are likely to have already been identified but not yet implemented. [31] A There were study results confirming that the combination of age at onset and tumor phenotype can provide an efficient model for identifying individuals with a high probability of carrying BRCA mutations and supporting the hypothesis that breast cancer in BRCA carriers is qualitatively distinct from other early-onset breast cancers and from late-onset, sporadic, breast carcinomas. [32] Tumors with mutations in the gene encoding the serine-threonine protein kinase BRAF are dependent on the MAPK signaling pathway for their growth, what offers an opportunity to test oncogene-targeted therapy. Mutations at the position V600 of BRAF were described in approximately 8% of all solid tumors, including 50% of melanomas, 30 to 70% of papillary thyroid carcinomas and 5 to 8% of colorectal adenocarcinomas. Specific BRAF kinase inhibitors are undergoing rapid clinical development and promising data on efficacy have been demonstrated in activated mutant BRAF V600 melanomas. This review article will address: (a) preclinical data on the antitumor activity of BRAF inhibitors in cell lines/ in vivo models and their opposing functions [33] The advent of targeted therapies has revolutionized the treatment of certain types of cancer. Identification of molecular targets on cancer cells has led to the design of novel drugs, which either used as single agents or in combination with chemotherapy, has prolonged survival in metastatic disease, or contributed to curative treatment in the adjuvant setting. A literature review was conducted to identify and present current knowledge on the molecular function of the HER2 receptor, its role in the pathogenesis of breast cancer and anti-HER2 targeted drugs in use or under development. Many molecular targets have been identified in breast cancer, with the HER family of receptors being the ones most extensively studied [34] There was published study results suggest that circulating plasma sKIT levels seem to function as a surrogate marker for TTP in gastrointestinal stromal tumor patients although additional studies are warranted to confirm and expand these findings[35] According to some study, it is reasonable to suggest that PSA density, velocity, doubling time and free to total PSA ratio or combining PSA with Gleason score shall greatly increase PSA specificity in detecting PC cases. Radioisotopic bone scan by SPET or PET can demonstrate osseous metastases at later stages of PC, but should also be applied in cases falsely considered as an early stage of PC, for better staging and for periodic follow-up of the disease. [36] The protein S-100B seems to be the best analyzed biomarker in melanoma and it has the potential to identify high-risk stage III melanoma patients who may benefit from adjuvant systematic treatment. Since an effective (adjuvant) therapy for loco-regional metastatic and disseminated melanoma is recently introduced, the use of S-100B seems to alter dramatically in the near future [37] Despite advances in genomics, proteomics and molecular pathology have generated many candidate biomarkers with potential clinical value, their use for cancer staging and personalization of therapy at the time of diagnosis could improve patient care. The challenge of translation from bench to bedside outside of the research setting has proved more difficult than might have been expected. Thus, understanding how and when
Biomarkers can be integrated into clinical care is crucial if to translate the promise into reality as reported already. [38] Biomarkers present the normal and/or disease state in humans. Genetic and epigenetic biomarkers assessed in easily accessible biological materials are useful in diagnosis, early onset or risk of developing cancer or to predict the treatment efficacy or clinical outcome of different human malignancies. Some of these markers are also expressed during early stages of the tumor development and hence provide an opportunity to develop intervention and treatment strategies. With the advances of research in clinical genetic and epigenetic, once validated, these markers can be utilized in clinical settings and to identify high risk populations through their implication in cancer diagnosis and risk assessment. [39] Tissue inhibitor of metalloproteinase-1 (TIMP-1) was evaluated in the pre-treatment serum of newly diagnosed patients with symptomatic myeloma and revealed that pre-treatment serum TIMP-1 is associated with advanced myeloma suggesting further evaluation of this molecule to better determine its prognostic potential in MM. [40] Significant number of metastases from hormone receptor-positive primary breast cancer do not respond to endocrine therapy. A study to assess how often hormone receptor status changes between primary and recurrent tumors and whether such a change might explain unresponsiveness to endocrine therapy indicated that loss of estrogen receptor (ER) expression in recurrent breast cancer should be considered as a cause for poor response to endocrine therapy in primarily ER-positive patients. [41] According to a review from thirteen categories of breast tumor markers considered, not all applications for these markers were supported due some demonstrated insufficient evidence to support routine use in clinical practice [42] The therapeutic strategy for breast cancer with the use of targeted drugs is, at present, mainly focused on coping with HER2 as currently for instance lapatinib and trastuzumab are in widespread use. [43] Treatment decisions for patients with lung cancer have historically been based on tumor histology. Mutations in EGFR and KRAS have been extensively studied, but the application of treatment according to the genetic potential of individual tumours may lead to substantial therapeutic improvements as stated in recent review. [44] The tissue concentrations of specific miRNAs have been associated with tumor invasiveness, metastatic potential, and other clinical characteristics for several types of cancers, including chronic lymphocytic leukemia, and breast, colorectal, hepatic, lung, pancreatic, and prostate cancers. The biologic roles of miRNAs in cancer also suggest a correlation with prognosis and therapeutic outcome indicating further investigation of these roles to lead to new approaches for the categorization, diagnosis, and treatment of human cancers. [45]

<table>
<thead>
<tr>
<th>Tumor Type</th>
<th>Biomarker</th>
</tr>
</thead>
<tbody>
<tr>
<td>Breast</td>
<td>ER/FR (estrogen receptor/progesterone receptor)</td>
</tr>
<tr>
<td></td>
<td>HER-2/neu</td>
</tr>
<tr>
<td></td>
<td>EGFR</td>
</tr>
<tr>
<td>Colorectal</td>
<td>KRAS</td>
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<tr>
<td></td>
<td>UGT1A1</td>
</tr>
<tr>
<td>Gastric</td>
<td>HER-2/neu</td>
</tr>
<tr>
<td></td>
<td>e-c-kit</td>
</tr>
<tr>
<td></td>
<td>CD20 Antigen</td>
</tr>
<tr>
<td></td>
<td>CD80</td>
</tr>
<tr>
<td></td>
<td>PIP111-PDGFRalpha</td>
</tr>
<tr>
<td>Lymphoma</td>
<td>Philadelphia Chromosome (BCR/ABL)</td>
</tr>
<tr>
<td></td>
<td>PML/RAR alpha</td>
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<tr>
<td></td>
<td>TP53T</td>
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<tr>
<td></td>
<td>UGT1A1</td>
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<tr>
<td></td>
<td>ALK</td>
</tr>
<tr>
<td>Lung</td>
<td>EGFR</td>
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<td></td>
<td>KRAS</td>
</tr>
<tr>
<td>Melanoma</td>
<td>BRAF</td>
</tr>
</tbody>
</table>

Table 2. Breif list of some molecular cancer biomarkers

Beyond the breif list above, other examples of biomarkers also include: Tumor suppressors lost in cancer (Examples: BRCA1, BRCA2), RNA (Examples: mRNA, microRNA), and proteins found in body fluids or tissue like prostate-specific antigen, and CA-125). [46]

VII. Conclusion

Today we have several biomarkers used in major fields and specialties of medicine. These are not only to help disease diagnosis and differential diagnosis, also to develop new drugs to target certain pathologies and to rationally use treatment measures available at the right time for the right patient. Keeping on utilizing the well established markers, we can also develop the newer and more sensitive and specific ones based on the present knowledge deepening supported by further clinical application and evidences collected from multicentere settings.
conclusion section is not required. Although a conclusion may review the main points of the paper, do not replicate the abstract as the conclusion. A conclusion might elaborate on the importance of the work or suggest applications and extensions.

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